# Boundary Controllability of the Korteweg-de Vries Equation on a Bounded Domain

## Eduardo Cerpa

Departamento de Matemática Universidad Técnica Federico Santa María Casilla 110-V, Valparaíso, Chile email: eduardo.cerpa@usm.cl

## Ivonne Rivas

Instituto de Matemática Pura e Aplicada (IMPA) Rio de Janeiro 22460-320, Brazil email: ivonriv@impa.br

## Bing-Yu Zhang

Department of Mathematical Sciences
University of Cincinnati
Cincinnati, Ohio 45221, USA
and
Yangtze Center of Mathematics
Sichuan University
Chengdu, China
email: zhangb@ucmail.uc.edu

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#### Abstract

This paper is devoted to study boundary controllability of the Korteweg-de Vries equation posed on a finite interval, in which, because of the third-order character of the equation, three boundary conditions are required to secure the well-posedness of the system. We consider the cases where one, two, or all three of those boundary data are employed as boundary control inputs. The system is first linearized around the origin and the corresponding linear system is shown to be exactly boundary controllable if using two or three boundary control inputs. In the case where only one control input is allowed to be used, the linearized system is known to be only *null* controllable if the single control input acts on the left end of the spatial domain. By contrast, if the single control input acts on the right end of the spatial domain, the linearized system is exactly controllable if and only if the length of the spatial domain does not belong to a set of *critical values*. Moreover, the nonlinear system is shown to be locally exactly boundary controllable via *contraction mapping principle* if the associated linearized system is exactly controllable.

**Key words.** Boundary control, Exact controllability, the Korteweg-de Vries equation, non-linear systems

AMS subject classifications. 93B05, 35Q53, 35Q53

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## 1 Introduction

In this paper we study a class of distributed parameter control systems described by the Korteweg-de Vries (KdV) equation posed on a finite domain with nonhomogeneous boundary conditions:

$$\begin{cases} y_t + y_x + y_{xxx} + yy_x = 0, & (x,t) \in (0,L) \times (0,T), \\ y(0,t) = h_1(t), y_x(L,t) = h_2(t), y_{xx}(L,t) = h_3(t), & t \in (0,T). \end{cases}$$
(1.1)

This system can be considered as a model for propagation of surface water waves in the situation where a wave-maker is putting energy in a finite-length channel from the left (x = 0) while the right end (x = L) of the channel is free (corresponding to the case  $h_2 = h_3 = 0$ ) (see [5]). Since the work of Colin and Ghidaglia in the late 1990's [5, 6, 7], the system (1.1) has been mainly studied for its well-posedness in the classical Sobolev space  $H^s(0, L)$  [13, 19]. So far, the system is known to be locally well-posed in the space  $H^s(0, L)$  for any  $s > -\frac{3}{4}$  as stated in the following theorem.

**Theorem A** [14]: Let  $s > -\frac{3}{4}$ , T > 0 and r > 0 be given with

$$s \neq \frac{2j-1}{2}, \quad j = 1, 2, 3, \cdots.$$

There exists  $T^* > 0$  such that for given s-compatible  $^1$  data

$$y_0 \in H^s(0,L), \quad h_1 \in H^{\frac{s+1}{3}}(0,T), \quad h_2 \in H^{\frac{s}{3}}(0,T), \quad h_3 \in H^{\frac{s-1}{3}}(0,T)$$

satisfying

$$||y_0||_{H^s(0,L)} + ||h_1||_{H^{\frac{s+1}{3}}(0,T)} + ||h_2||_{H^{\frac{s}{3}}(0,T)} + ||h_3||_{H^{\frac{s-1}{3}}(0,T)} \le r,$$

then (1.1) admits a unique solution

$$y \in C([0,T^*];H^s(0,L)) \cap L^2(0,T^*;H^{s+1}(0,L))$$

<sup>&</sup>lt;sup>1</sup>The reader is referred to [14] for the precise definition of s-compatibility for the IBVP (1.1).

satisfying the initial condition

$$y|_{t=0} = y_0.$$

Moreover, the solution y depends Lipschitz continuously on  $y_0$  and  $h_j$ , j = 1, 2, 3 in the corresponding spaces.

In this paper we are interested in studying the IBVP (1.1) from a control point of view:

How solutions of the system (1.1) can be influenced by choosing appropriate control inputs  $h_i$ , j = 1, 2, 3?

In particular, we are concerned with the following exact boundary control problem.

**Exact Control Problem:** Given T > 0 and  $y_0, y_T \in L^2(0, L)$ , can one find appropriate control inputs  $h_j$ , j = 1, 2, 3 such that the corresponding solution y of (1.1) satisfies

$$y|_{t=0} = y_0, \qquad y|_{t=T} = y_T$$
?

Boundary control problems for the KdV equation on a finite domain have been extensively studied in the past (see [25, 20, 26, 21, 8, 3, 9, 4, 23] and the references therein). Most of those works have been focused on the following system

$$\begin{cases} u_t + u_x + u_{xxx} + uu_x = 0, & (x,t) \in (0,L) \times (0,T), \\ u(0,t) = g_1(t), & u(L,t) = g_2(t), & u_x(L,t) = g_3(t), & t \in (0,T) \end{cases}$$

$$(1.2)$$

which possess a different set of boundary conditions than those of the system (1.1). Controllability of this system was first studied by Rosier [20] in 1997 using only one control input  $g_3$ :

$$\begin{cases}
 u_t + u_x + u_{xxx} + uu_x = 0, & (x,t) \in (0,L) \times (0,T), \\
 u(0,t) = 0, & u(L,t) = 0, & u_x(L,t) = g_3(t), & t \in (0,T).
\end{cases}$$
(1.3)

It was discovered rather surprisingly that whether the associated linear system

$$\begin{cases} u_t + u_x + u_{xxx} = 0, & (x,t) \in (0,L) \times (0,T), \\ u(0,t) = 0, & u(L,t) = 0, & u_x(L,t) = g_3(t), & t \in (0,T), \end{cases}$$
(1.4)

is exactly controllable depends on the length L of the spatial domain (0, L). More precisely, Rosier [20] showed that the linear system is exactly controllable in the space  $L^2(0, L)$  if and only if

$$L \notin \mathcal{S} := \left\{ 2\pi \sqrt{\frac{k^2 + kl + l^2}{3}}; \, k, l \in \mathbb{N}^* \right\}. \tag{1.5}$$

With the linear result in hand and using the contraction mapping principle, Rosier [20] showed further that the nonlinear system (1.3) is locally exactly controllable in the space  $L^2(0, L)$  so long as  $L \notin \mathcal{S}$ .

**Theorem B** (Rosier [20]): Let T > 0 be given and assume  $L \notin S$ . There exists r > 0 such that for any  $u_0$ ,  $u_T \in L^2(0, L)$  with

$$||u_0||_{L^2(0,L)} + ||u_T||_{L^2(0,L)} \le r,$$

there exists  $g_3 \in L^2(0,T)$  such that the system (1.3) admits a unique solution

$$u \in C([0,T]; L^2(0,L)) \cap L^2(0,T; H^1(0,L))$$

satisfying

$$u|_{t=0} = u_0, \qquad u|_{t=T} = u_T.$$

The system (1.2) was later studied by Glass and Guerrero [10] for its boundary controllability using only  $g_2$  as a control input.

$$\begin{cases}
 u_t + u_x + u_{xxx} + uu_x = 0, & (x,t) \in (0,L) \times (0,T), \\
 u(0,t) = 0, & u(L,t) = g_2(t), & u_x(L,t) = 0, & t \in (0,T).
\end{cases}$$
(1.6)

They showed that the corresponding linear system

$$\begin{cases} u_t + u_x + u_{xxx} = 0, & (x,t) \in (0,L) \times (0,T), \\ u(0,t) = 0, & u(L,t) = g_2(t), & u_x(L,t) = 0, \quad t \in (0,T), \end{cases}$$
(1.7)

is exactly controllable in the space  $L^2(0,L)$  if and only if  $L \notin \mathcal{N}$  where

$$\mathcal{N} = \left\{ L \in \mathbb{R}^+ : L^2 = -(a^2 + ab + b^2) \text{ with } a, b \in \mathbb{C} \right.$$

$$\text{satisfying} \quad ae^a = be^b = -(a+b)e^{-(a+b)} \right\}. \tag{1.8}$$

Then the nonlinear system (1.6) was shown to be locally exactly controllable in the space  $L^2(0, L)$  if  $L \notin \mathcal{N}$ .

**Theorem C** (Glass and Guerrero [10]) Let T > 0 and  $L \notin \mathcal{N}$  be given. There exists r > 0 such that for any  $u_0$ ,  $u_T \in L^2(0, L)$  with

$$||u_0||_{L^2(0,L)} + ||u_T||_{L^2(0,L)} \le r,$$

one can find  $g_2 \in H^{\frac{1}{6}-\epsilon}(0,T)$  for any  $\epsilon > 0$  such that the system (1.6) admits a solution

$$u \in C([0,T]; L^2(0,L)) \cap L^2(0,T; H^1(0,L))$$

satisfying

$$u|_{t=0} = u_0, \qquad u|_{t=T} = u_T.$$

While the critical length phenomenon occurs when a single control input (either  $g_2$  or  $g_3$ ) is used, it will not happen, however, if more than one control inputs are allowed to be used. It was already pointed out in [20] that the linear system associated to (1.2) is exactly controllable for any L > 0 if both  $g_2$  and  $g_3$  are allowed to be used as control inputs. Moreover, the nonlinear system

$$\begin{cases}
 u_t + u_x + u_{xx} + uu_x = 0, & (x,t) \in (0,L) \times (0,T), \\
 u(0,t) = 0, & u(L,t) = g_2(t), & u_x(L,t) = g_3(t), & t \in (0,T),
\end{cases}$$
(1.9)

was shown in [20] to be locally exactly controllable in  $L^2(0, L)$ .

**Theorem D** (Rosier [20]): Let T > 0 and L > 0 be given and k > 0 be an integer. There exists  $\delta > 0$  such that for any  $u_0, u_T \in L^2(0, L)$  with

$$||u_0||_{L^2(0,L)} + ||u_T||_{L^2(0,L)} \le \delta,$$

there exist  $g_2 \in H_0^k(0,T)$  and  $g_3 \in L^2(0,T)$  such that the system (1.9) admits a solution

$$u \in C([0,T]; L^2(0,L)) \cap L^2(0,T; H^1(0,L))$$

satisfying

$$u|_{t=0} = u_0, \qquad u|_{t=T} = u_T.$$

In [10] Glass and Guerrero considered the system (1.2) using  $g_1$  and  $g_2$  as control inputs.

$$\begin{cases}
 u_t + u_x + u_{xxx} + uu_x = 0, & (x,t) \in (0,L) \times (0,T), \\
 u(0,t) = g_1, & u(L,t) = g_2(t), & u_x(L,t) = 0, & t \in (0,T),
\end{cases}$$
(1.10)

and showed that the system is also locally exactly controllable for any L > 0.

**Theorem E** (Glass and Guerrero [9]) Let T > 0 and L > 0 be given. There exists r > 0 such that for any  $u_0$ ,  $u_T \in L^2(0, L)$  with

$$||u_0||_{L^2(0,L)} + ||u_T||_{L^2(0,L)} \le r,$$

there exist  $g_1$  and  $g_2 \in L^2(0,T)$  such that the system (1.10) admits a solution

$$u \in C([0,T]; H^{-1}(0,L)) \cap L^2(0,T; L^2(0,L))$$

satisfying

$$u|_{t=0} = u_0, \qquad u|_{t=T} = u_T.$$

Another interesting controllability result regarding the system (1.2) is that it is only *null controllable* if the control acts from the left side of the spatial domain (0, L), which was proved by Rosier [21], Glass and Guerrero [9].

**Theorem F**([21, 9]) Let T > 0 and L > 0 be given. Let

$$v \in C([0,T]; L^2(0,L)) \cap L^2(0,T; H^1(0,L))$$

satisfy

$$\begin{cases} v_t + v_x + vv_x + v_{xxx} = 0, \ v(x,0) = v_0(x), & (x,t) \in (0,L) \times (0,T), \\ v(0,t) = v(L,t) = v_x(L,t) = 0, & t \in (0,T). \end{cases}$$
 (1.11)

Then, there exists  $\delta > 0$  such that for any  $u_0 \in L^2(0,L)$  with

$$||u_0 - v_0||_{L^2(0,L)} \le \delta,$$

there exists  $g_1 \in H^{\frac{1}{2}-\epsilon}(0,T)$  for any  $\epsilon > 0$  such that the system (1.2) admits a solution

$$u \in C([0,T]; L^2(0,L)) \cap L^2(0,T; H^1(0,L))$$

satisfying

$$u|_{t=0} = u_0, \qquad u|_{t=T} = v|_{t=T}.$$

While boundary controllability of the system (1.2) has been well studied, there is very few results for the system (1.1). To our knowledge, the only known result is due to Guilleron [11]. He has considered the system (1.1) with  $h_2 = h_3 = 0$  and shown the corresponding linear system

$$\begin{cases} y_t + y_x + y_{xxx} = 0, & (x,t) \in (0,L) \times (0,T), \\ y(0,t) = h_1(t), & y_x(L,t) = y_{xx}(L,t) = 0 \end{cases}$$

is null controllable by applying a Carleman estimates approach to obtain the needed observability inequality.

The purpose of this paper is to fill the gap and to determine if the system (1.1) possesses controllability results similar to those established for system (1.2). Naturally one would like to try the same approaches that have worked effectively for system (1.2). However, one will encounter some difficulties that demand special attention and some new tools will be needed. In particular, when we use only  $h_2$  as a control input, the linear system associated to (1.1) is

$$\begin{cases} y_t + y_x + y_{xxx} = 0, & y(x,0) = y_0(x), & (x,t) \in (0,L) \times (0,T), \\ y(0,t) = 0, & y_x(L,t) = h_2(t), & y_{xx}(L,t) = 0, & t \in (0,T) \end{cases}$$
(1.12)

and its adjoint system is given by

$$\begin{cases} \varphi_t + \varphi_x + \varphi_{xxx} = 0, & \varphi(x, T) = \varphi_T(x), \quad (x, t) \in (0, L) \times (0, T), \\ \varphi(L, t) + \varphi_{xx}(L, t) = 0, & \varphi_x(0, t) = 0, & \varphi(0, t) = 0, \quad t \in (0, T). \end{cases}$$
(1.13)

It is well-known that the exact controllability of system (1.12) is equivalent to the following observability inequality for the adjoint system (1.13)

$$\|\varphi_T\|_{L^2(0,L)} \le C\|\varphi_x(L,\cdot)\|_{L^2(0,T)}.$$
(1.14)

However, the usual multiplier method and compactness arguments as those used in dealing with the control of system (1.4) only lead to

$$\|\varphi_T\|_{L^2(0,L)} \le C_1 \|\varphi_x(L,\cdot)\|_{L^2(0,T)} + C_2 \|\varphi(L,\cdot)\|_{L^2(0,T)}. \tag{1.15}$$

How to remove the extra term in (1.15) presents a challenge and demands a new tool. This new tool turns out to be the hidden regularity (or the sharp Kato smoothing property [12]) for solutions of the KdV equation. Specially, as we will demonstrate later in this paper, for solutions of the system (1.13), the following inequality holds

$$\sup_{0 < x < L} \|\varphi(x, \cdot)\|_{H^{\frac{1}{3}}(0, T)} \le C \|\varphi_T\|_{L^2(0, L)}, \tag{1.16}$$

which will play a crucial role in validating the observability estimate (1.14).

In this paper, we will first consider case that only  $h_2$  is employed as a control input and show that the system

$$\begin{cases} y_t + y_x + y_{xxx} + yy_x = 0, (x,t) \in (0,L) \times (0,T), \\ y(0,t) = 0, \ y_x(L,t) = h_2(t), \ y_{xx}(L,t) = 0, \quad t \in (0,T) \end{cases}$$
(1.17)

is locally exactly controllable as long as  $L \notin \mathbb{F}$  where

$$\mathbb{F} = \left\{ L \in \mathbb{R}^+ : L^2 = -(a^2 + ab + b^2) \text{ with } a, b \in \mathbb{C} \text{ satisfying } \frac{e^a}{a^2} = \frac{e^b}{b^2} = \frac{e^{-(a+b)}}{(a+b)^2} \right\}. \quad (1.18)$$

**Theorem 1.1.** Let T > 0 and  $L \notin \mathbb{F}$  be given. There exists  $\delta > 0$  such that for any  $y_0, y_T \in L^2(0,L)$  with

$$||y_0||_{L^2(0,L)} + ||y_T||_{L^2(0,L)} < \delta,$$

one can find  $h_2 \in L^2(0,T)$  such that the system (1.17) admits a unique solution

$$y \in C([0,T]; L^2(0,L)) \cap L^2(0,T; H^1(0,L))$$

satisfying

$$y|_{t=0} = y_0, y|_{t=T} = y_T.$$

**Remark 1.2.** It can be proven that the set  $\mathbb{F}$  is nonempty and countable. The proof follows [10] where the authors proved that the set  $\mathcal{N}$  (see (1.8)) is nonempty and countable.

Instead of employing control input  $h_2$ , one can just use the control input  $h_3$ :

$$\begin{cases} y_t + y_x + y_{xxx} + yy_x = 0, (x,t) \in (0,L) \times (0,T), \\ y(0,t) = 0, y_x(L,t) = 0, y_{xx}(L,t) = h_3(t), \quad t \in (0,T). \end{cases}$$
(1.19)

The corresponding system is also locally exactly controllable if the length L of the interval (0, L) does not belong to the set  $\mathcal{N}$  as defined in (1.8).

**Theorem 1.3.** Let T > 0 and  $L \notin \mathcal{N}$  be given. There exists  $\delta > 0$  such that for any  $y_0, y_T \in L^2(0,L)$  with

$$||y_0||_{L^2(0,L)} + ||y_T||_{L^2(0,L)} \le \delta,$$

there exists  $h_3 \in H^{-\frac{1}{3}}(0,T)$  such that the system (1.19) admits a unique solution

$$y \in C([0,T]; L^2(0,L)) \cap L^2(0,T; H^1(0,L))$$

satisfying

$$y|_{t=0} = y_0, \qquad y|_{t=T} = y_T.$$

Similar to the system (1.2), the critical length phenomenon will not occur if more than one control inputs are employed. For the system where  $h_1$  and  $h_2$  are used as control inputs,

$$\begin{cases} y_t + y_x + y_{xxx} + yy_x = 0, (x,t) \in (0,L) \times (0,T), \\ y(0,t) = h_1(t), y_x(L,t) = h_2(t), y_{xx}(L,t) = 0, \quad t \in (0,T) \end{cases}$$
(1.20)

we have the following local exact controllability result.

**Theorem 1.4.** Let T > 0 and L > 0 be given. There exists  $\delta > 0$  such that for any  $y_0, y_T \in L^2(0,L)$  with

$$||y_0||_{L^2(0,L)} + ||y_T||_{L^2(0,L)} \le \delta,$$

one can find  $h_1 \in H^{\frac{1}{3}}(0,T)$  and  $h_2 \in L^2(0,T)$  such that the system (1.20) admits a unique solution

$$y \in C([0,T]; L^2(0,L)) \cap L^2(0,T; H^1(0,L))$$

satisfying

$$y|_{t=0} = y_0, \qquad y|_{t=T} = y_T.$$

For the system using both  $h_2$  and  $h_3$  as control inputs,

$$\begin{cases} y_t + y_x + y_{xxx} + yy_x = 0, (x,t) \in (0,L) \times (0,T), \\ y(0,t) = 0, y_x(L,t) = h_2(t), y_{xx}(L,t) = h_3(t), \quad t \in (0,T), \end{cases}$$
(1.21)

we have the following local exact controllability result.

**Theorem 1.5.** Let T > 0 and L > 0 be given. There exists  $\delta > 0$  such that for any  $y_0, y_T \in L^2(0,L)$  with

$$||y_0||_{L^2(0,L)} + ||y_T||_{L^2(0,L)} \le \delta,$$

one can find  $h_3 \in H^{-\frac{1}{3}}(0,T)$  and  $h_2 \in L^2(0,T)$  such that the system (1.21) admits a unique solution

$$y\in C([0,T];L^2(0,L))\cap L^2(0,T;H^1(0,L))$$

satisfying

$$y|_{t=0} = y_0, \qquad y|_{t=T} = y_T.$$

For the system using  $h_1$  and  $h_3$  as control inputs,

$$\begin{cases} y_t + y_x + y_{xxx} + yy_x = 0, (x,t) \in (0,L) \times (0,T), \\ y(0,t) = h_1(t), \ y_x(L,t) = 0, \ y_{xx}(L,t) = h_3(t), \quad t \in (0,T) \end{cases}$$
(1.22)

we have

**Theorem 1.6.** Let T > 0 and L > 0 be given. There exists  $\delta > 0$  such that for any  $y_0, y_T \in L^2(0,L)$  with

$$||y_0||_{L^2(0,L)} + ||y_T||_{L^2(0,L)} \le \delta,$$

one can find  $h_1 \in H^{\frac{1}{3}}(0,T)$  and  $h_3 \in H^{-\frac{1}{3}}(0,T)$  such that the system (1.22) admits a unique solution

$$y\in C([0,T];L^2(0,L))\cap L^2(0,T;H^1(0,L))$$

satisfying

$$y|_{t=0} = y_0, y|_{t=T} = y_T.$$

If all three boundary control inputs are allowed to be used, then we can show that system (1.1) is locally exactly controllable around any smooth solution of the KdV equation.

**Theorem 1.7.** Let T>0 and L>0 be given. Assume that  $u\in C^{\infty}(\mathbb{R},H^{\infty}(\mathbb{R}))$  satisfies

$$u_t + u_x + uu_x + u_{xxx} = 0, \quad (x, t) \in \mathbb{R} \times \mathbb{R}.$$

Then there exists  $\delta > 0$  such that for any  $y_0, y_T \in L^2(0, L)$ , satisfying

$$||y_0 - u(\cdot, 0)||_{L^2(0,L)} + ||y_T - u(\cdot, T)||_{L^2(0,L)} \le \delta$$

one can find control inputs

$$h_1 \in H^{\frac{1}{3}}(0,T), \quad h_2 \in L^2(0,T), \quad h_3 \in H^{-\frac{1}{3}}(0,T)$$

such that (1.1) admits a unique solution

$$y \in C([0,T];L^2(0,L)) \cap L^2(0,T;H^1(0,L))$$

satisfying

$$y|_{t=0} = y_0, \qquad y|_{t=T} = y_T.$$

Finally, with the help of some hidden regularity properties for solutions of the KdV equation we can improve some controllability results for the system (1.2). Concerning system

$$\begin{cases} u_t + u_x + u_{xxx} + uu_x = 0, & (x,t) \in (0,L) \times (0,T), \\ u(0,t) = 0, & u(L,t) = g_2(t), & u_x(L,t) = 0, & t \in (0,T) \end{cases}$$
(1.23)

we can show that the control input  $g_2$ , used in Theorem C, belongs in fact to the space  $H^{\frac{1}{3}}(0,T)$ .

**Theorem 1.8.** Let T > 0 and  $L \notin \mathcal{N}$  be given. There exists r > 0 such that for any  $u_0, u_T \in L^2(0,L)$  with

$$||u_0||_{L^2(0,L)} + ||u_T||_{L^2(0,L)} \le r,$$

there exists  $g_2 \in H^{\frac{1}{3}}(0,T)$  such that the system (1.23) admits a solution

$$u \in C([0,T]; L^2(0,L)) \cap L^2(0,T; H^1(0,L))$$

satisfying

$$u|_{t=0} = u_0, \qquad u|_{t=T} = u_T.$$

Regarding system

$$\begin{cases}
 u_t + u_x + u_{xxx} + uu_x = 0, & (x,t) \in (0,L) \times (0,T), \\
 u(0,t) = 0, & u(L,t) = g_2(t), & u_x(L,t) = g_3(t), & t \in (0,T),
\end{cases}$$
(1.24)

the Theorem D can be improved as follows.

**Theorem 1.9.** Let T > 0 and L > 0 be given. There exists  $\delta > 0$  such that for any  $u_0$ ,  $u_T \in L^2(0,L)$  with

$$||u_0||_{L^2(0,L)} + ||u_T||_{L^2(0,L)} \le \delta,$$

there exist  $g_2 \in H^{\frac{1}{3}}(0,T)$  and  $g_3 \in L^2(0,T)$  such that system (1.24) admits a solution

$$u \in C([0,T]; L^2(0,L)) \cap L^2(0,T; H^1(0,L))$$

satisfying

$$u|_{t=0} = u_0, \qquad u|_{t=T} = u_T.$$

Moreover, we also consider system (1.2) using only two control inputs  $g_1$  and  $g_2$ 

$$\begin{cases}
 u_t + u_x + u_{xxx} + uu_x = 0, & (x,t) \in (0,L) \times (0,T), \\
 u(0,t) = g_1(t), & u(L,t) = g_2(t), & u_x(L,t) = 0, & t \in (0,T)
\end{cases}$$
(1.25)

which has not been studied in the literature before. We can show that the critical length phenomenon will not occur for this system neither.

**Theorem 1.10.** Let T > 0 and L > 0 be given. There exists  $\delta > 0$  such that for any  $u_0$ ,  $u_T \in L^2(0,L)$  with

$$||u_0||_{L^2(0,L)} + ||u_T||_{L^2(0,L)} \le \delta,$$

one can find  $g_1 \in H^{\frac{1}{3}}(0,T)$  and  $g_2 \in H^{\frac{1}{3}}(0,T)$  such that the system (1.25) admits a solution

$$u \in C([0,T]; L^2(0,L)) \cap L^2(0,T; H^1(0,L))$$

satisfying

$$u|_{t=0} = u_0, \qquad u|_{t=T} = u_T.$$

The paper is organized as follows. In Section 2, we present various linear estimates including hidden regularities for solutions of the linear systems associated to (1.1) and (1.2) which will play important roles in establishing our exact controllability results in this paper. The associated linear systems are shown to be exactly controllable in Section 3 while the nonlinear systems are shown to be locally exactly controllable using the standard contraction mapping principle in Section 4. Finally, in Section 5 we provide some conclusion remarks together with some open problems for further studies. The paper is ended with an appendix where the proofs of some technical lemmas used in the paper are furnished.

## 2 Linear estimates

### 2.1 The forward linear system

In this subsection, we consider the following linear problem associated to the nonlinear system (1.1)

$$\begin{cases} y_t + y_x + y_{xxx} = f, & y(x,0) = y_0(x), & x \in (0,L), \ t > 0, \\ y(0,t) = h_1(t), \ y_x(L,t) = h_2(t), \ y_{xx}(L,t) = h_3(t), & t > 0. \end{cases}$$
(2.1)

In the case  $h_1 = h_2 = h_3 = 0$  and f = 0, the solution y can be written as

$$y = W_0(t)y_0$$
.

Here  $W_0(t)$  is the  $C_0$ -semigroup in the space  $L^2(0,L)$  (see [16]) generated by the linear operator

$$A\psi = -\psi''' - \psi'$$

whose domain is

$$\mathcal{D}(A) = \{ \psi \in H^3(0, L) : \psi(0) = \psi'(L) = \psi''(L) = 0 \}.$$

The solution y of (2.1), when  $h_1 = h_2 = h_3 = 0$  and  $y_0 = 0$ , is given by

$$y = \int_0^t W_0(t - \tau) f(\tau) d\tau$$

while in the case  $y_0 = 0$  and f = 0, it has the form

$$y = W_{bdr}(t)\vec{h}$$

where  $\vec{h} = (h_1, h_2, h_3)$  and  $W_{bdr}(t)$  is the associated boundary integral operator defined in [13, 14].

As it has been demonstrated in [13, 14] the linear system (2.1) is well-posed in the space  $H^s(0, L)$  for any  $0 \le s \le 3$  with

$$y_0 \in H^s(0,L), f \in W^{\frac{s}{3},1}(0,T;L^2(0,L))$$

and

$$\vec{h} = (h_1, h_2, h_3) \in \mathcal{H}^s_{loc}(\mathbb{R}^+) := H^{\frac{s+1}{3}}_{loc}(\mathbb{R}^+) \times H^{\frac{s}{3}}_{loc}(\mathbb{R}^+) \times H^{\frac{s-1}{3}}_{loc}(\mathbb{R}^+).$$

In particular, in the case s = 0, the result can be stated as follows.

**Proposition 2.1.** Let T > 0 be given, for any  $y_0 \in L^2(0,L)$ ,  $f \in L^1(0,T;L^2(0,L))$  and

$$(h_1, h_2, h_3) \in \mathcal{H}_T := H^{\frac{1}{3}}(0, T) \times L^2(0, T) \times H^{-\frac{1}{3}}(0, T),$$

the IBVP (2.1) admits a unique solution

$$y \in X_T := C([0,T]; L^2(0,L)) \cap L^2(0,T; H^1(0,L)).$$

Moreover, there exists C > 0 such that

$$||y||_{X_T} \le C \left( ||f||_{L^1(0,T;L^2(0,L))} + ||y_0||_{L^2(0,L)} + ||(h_1,h_2,h_3)||_{\mathcal{H}_T} \right).$$

In addition, the solution y of (2.1) possesses the following hidden (or sharp trace) regularities.

**Proposition 2.2.** Let T > 0 be given. For any  $y_0 \in L^2(0,L)$ ,  $f \in L^1(0,T;L^2(0,L))$  and  $(h_1,h_2,h_3) \in \mathcal{H}_T$ , the solution y of the system (2.1) satisfies

$$\sup_{0 < x < L} \|\partial_x^j y(x, \cdot)\|_{H^{\frac{1-j}{3}}(0,T)} \le C_j \left( \|f\|_{L^1(0,T;L^2(0,L))} + \|y_0\|_{L^2(0,L)} + \|(h_1, h_2, h_3)\|_{\mathcal{H}_T} \right)$$
(2.2)

for j = 0, 1, 2.

Next proposition states similar hidden (or sharp trace) regularity results for the linear system

$$\begin{cases}
 u_t + u_x + u_{xxx} = f, & u(x,0) = u_0(x), & x \in (0,L), \ t > 0, \\
 u(0,t) = g_1(t), \ u(L,t) = g_2(t), \ u_x(L,t) = g_3(t), & t > 0,
\end{cases}$$
(2.3)

associated to (1.2).

**Proposition 2.3.** Let T>0 be given, for any  $u_0 \in L^2(0,L)$ ,  $f \in L^1(0,T;L^2(0,L))$  and

$$(g_1, g_2, g_3) \in \mathcal{G}_T := H^{\frac{1}{3}}(0, T) \times H^{\frac{1}{3}}(0, T) \times L^2(0, T),$$

the IBVP (2.3) admits a unique solution  $u \in X_T$ . Moreover, there exists C > 0 such that

$$||u||_{X_T} \le C \left( ||f||_{L^1(0,T;L^2(0,L))} + ||y_0||_{L^2(0,L)} + ||(g_1,g_2,g_3)||_{\mathcal{G}_T} \right).$$

In addition, the solution u possesses the following sharp trace estimates

$$\sup_{0 < x < L} \|\partial_x^j u(x, \cdot)\|_{H^{\frac{1-j}{3}}(0,T)} \le C_j \left( \|f\|_{L^1(0,T;L^2(0,L))} + \|u_0\|_{L^2(0,L)} + \|(g_1, g_2, g_3)\|_{\mathcal{G}_T} \right)$$
(2.4)

for j = 0, 1, 2.

The proofs of Proposition 2.2 and Proposition 2.3 can be found in [28] (cf. also [1, 2, 14]).

**Remark 2.4.** Systems (2.1) and (2.3) are equivalent in the following sense: for given  $\{y_0, f, h_1, h_2, h_3\}$  one can find  $\{u_0, f, g_1, g_2, g_3\}$  such that the corresponding solution y of (2.1) is exactly the same as the corresponding u for system (2.3) and vice versa. Indeed, for given  $y_0 \in L^2(0, L)$ ,  $f \in L^1(0,T;L^2(0,L))$  and  $\vec{h} \in \mathcal{H}_T$ , system (2.1) admits a unique solution  $y \in X_T$ . Let  $u_0 = y_0$ , and set

$$g_1(t) = h_1(t), \quad g_2(t) = y(L, t), \quad g_3(t) = h_2(t).$$

Then, according to (2.2), we have  $\vec{g} \in \mathcal{G}_T$ . Because of the uniqueness of the IBVP (2.3), with such selected  $(u_0, f, g_1, g_2, g_3)$ , the corresponding solution  $u \in X_T$  of (2.3) must be equal to y since y also solves (2.3) with the given auxiliary data  $(u_0, f, g_1, g_2, g_3)$ . On the other hand, for any given  $u_0 \in L^2(0, L)$ ,  $f \in L^1(0, T; L^2(0, L))$  and  $\vec{g} \in \mathcal{G}_T$ , let  $u \in X_T$  be the corresponding solution of the system (2.3). By (2.4), we have  $u_{xx}(L, \cdot) \in H^{-\frac{1}{3}}(0, T)$ . Thus, if set  $y_0 = u_0$  and

$$h_1(t) = g_1(t), \quad h_2(t) = g_3(t), \quad h_3(t) = u_{xx}(L, t),$$

then  $\vec{h} \in \mathcal{H}_T$  and the corresponding solution  $y \in X_T$  of (2.1) must be equal to u which also solves (2.1) with the auxiliary data  $(y_0, f, \vec{h})$ .

## 2.2 The backward adjoint linear system

In this subsection, we consider the backward adjoint system of (2.1)

$$\begin{cases} \psi_t + \psi_x + \psi_{xxx} = 0, & \psi(x, T) = \psi_T(x), & (x, t) \in (0, L) \times (0, T), \\ \psi(0, t) = 0, & \psi_x(0, t) = 0, & \psi(L, t) + \psi_{xx}(L, t) = 0, & t \in (0, T), \end{cases}$$
(2.5)

which (by transformation x' = L - x, t' = T - t) is equivalent to the following forward system

$$\begin{cases} \varphi_t + \varphi_x + \varphi_{xxx} = 0, & \varphi(x,0) = \varphi_0(x), & (x,t) \in (0,L) \times (0,T), \\ \varphi(L,t) = 0, & \varphi_x(L,t) = 0, & \varphi(0,t) + \varphi_{xx}(0,t) = 0, & t \in (0,T). \end{cases}$$
(2.6)

The solution of (2.6) can be written as

$$\varphi(x,t) = S(t)\varphi_0$$

where S(t) is the  $C_0$  semigroup in the space  $L^2(0,L)$  generated by the operator

$$A_1 f = -f' - f'''$$

with the domain

$$\mathcal{D}(A_1) = \{ f \in H^3(0, L) : f(0) + f''(0) = 0, f(L) = f'(L) = 0 \}.$$

**Proposition 2.5.** For any  $\varphi_0 \in L^2(0,L)$  the IBVP (2.6) admits a unique solution  $\varphi \in X_T$ . Moreover, there exists C > 0 such that

$$\|\varphi\|_{X_T} \le C \|\varphi_0\|_{L^2(0,L)}$$

and

$$\int_0^T \left( |\varphi(0,t)|^2 + |\varphi_x(0,t)|^2 \right) dt \le C \|\varphi_0\|_{L^2(0,L)}^2.$$

**Proof.** The proof is very similar to that of [20] and is therefore omitted.

Thus, when  $\varphi_0 \in L^2(0, L)$ , the corresponding solution  $\varphi$  has the trace  $\varphi(L, \cdot) \in L^2(0, T)$ . The next theorem reveals that  $\varphi$  has a stronger trace regularity:  $\varphi(L, \cdot) \in H^{\frac{1}{3}}(0, T)$ . It will play an important role to establish exact controllability of the system (1.1) as shown in the next section.

**Theorem 2.6** (Hidden regularities). For any  $\varphi_0 \in L^2(0,L)$ , the solution  $\varphi \in X_T$  of IBVP (2.6) possesses the following sharp trace properties

$$\sup_{0 \le x \le L} \|\partial_x^j \varphi(x, \cdot)\|_{H^{\frac{1-j}{3}}(0, T)} \le C_j \|\varphi_0\|_{L^2(0, L)}$$

for j = 0, 1, 2.

**Remark 2.7.** Equivalently the solutions of the system (2.3) has the following sharp trace estimates:

$$\sup_{0 \le x \le L} \|\partial_x^j \psi(x, \cdot)\|_{H^{\frac{1-j}{3}}(0, T)} \le C_j \|\psi_T\|_{L^2(0, L)}$$

for j = 0, 1, 2.

To prove Theorem 2.6, we first consider the following linear system

$$\begin{cases} w_t + w_{xxx} = f, & (x,t) \in (0,L) \times (0,T), \\ w_{xx}(0,t) = k_1(t), & w(L,t) = k_2(t), & w_x(L,t) = k_3(t), & t \in (0,T), \\ w(x,0) = w_0(x), & x \in (0,L). \end{cases}$$
(2.7)

**Proposition 2.8.** If  $w_0 \in L^2(0,L)$ ,  $f \in L^1(0,T;L^2(0,L))$  and  $\vec{k} := (k_1,k_2,k_3) \in \mathcal{K}_T$  with  $\mathcal{K}_T = H^{-\frac{1}{3}}(0,T) \times H^{\frac{1}{3}}(0,T) \times L^2(0,T)$ , then the system (2.7) admits a unique solution  $w \in X_T$  which, in addition, has the hidden (or sharp trace) regularities

$$\partial_x^j w \in L_x^{\infty}(0, L, H^{\frac{1-j}{3}}(0, T)) \text{ for } j = 0, 1, 2.$$

Moreover, there exist constants C > 0,  $C_j > 0$ , j = 0, 1, 2 such that

$$||w||_{X_T} \le C \left( ||w_0||_{L^2(0,L)} + ||\vec{k}||_{\mathcal{K}_T} + ||f||_{L^1(0,T;L^2(0,L))} \right),$$

where,

$$\|\vec{k}\|_{\mathcal{K}_T} := \left( \|k_1\|_{H^{-\frac{1}{3}}(0,T)}^2 + \|k_2\|_{H^{\frac{1}{3}}(0,T)}^2 + \|k_3\|_{L^2(0,T)}^2 \right)^{1/2}$$

and

$$\sup_{0 < x < L} \|\partial_x^j w(x, \cdot)\|_{H^{\frac{1-j}{3}}(0, T)} \le C_j \left( \|w_0\|_{L^2(0, L)} + \|\vec{k}\|_{\mathcal{K}_T} + \|f\|_{L^1(0, T; L^2(0, L))} \right)$$

for j = 0, 1, 2.

**Proof.** The proof is similar to the one in [28]. Its sketch will be presented in the Appendix for the convenience of the interested readers.

Now we turn to prove Theorem 2.6.

#### Proof of Theorem 2.6. Let

$$\mathcal{X}_T := \left\{ u \in X_T; \quad \partial_x^j u \in L_x^{\infty}(0, L; H^{\frac{1-j}{3}}(0, T)), \ j = 0, 1, 2 \right\}$$

which is a Banach space equipped with the norm

$$||u||_{\mathcal{X}_T} := ||u||_{X_T} + \sum_{j=0}^2 ||\partial_x^j u||_{L_x^{\infty}(0,L;H^{\frac{1-j}{3}}(0,T))}.$$

According to Proposition 2.8, for any  $v \in \mathcal{X}_{\beta}$  where  $0 < \beta \leq T$ , and any  $\varphi_0 \in L^2(0, L)$ , the system

$$\begin{cases} w_t + w_{xxx} = -v_x, & (x,t) \in (0,L) \times (0,\beta), \\ w_{xx}(0,t) = -v(0,t), & w(L,t) = 0, \quad w_x(L,t) = 0, \quad t \in (0,\beta), \\ w(x,0) = \varphi_0(x), & x \in (0,L) \end{cases}$$
 (2.8)

admits a unique solution  $w \in \mathcal{X}_{\beta}$  and, moreover,

$$||w||_{\mathcal{X}_{\beta}} \le C \left( ||\varphi_0||_{L^2(0,L)} + ||v(0,\cdot)||_{H^{-\frac{1}{3}}(0,\beta)} + ||v_x||_{L^1(0,\beta;L^2(0,L))} \right)$$

where the constant C depends only on T. As we have

$$||v_x||_{L^1(0,\beta;L^2(0,L))} \le \beta^{\frac{1}{2}} ||v||_{\mathcal{X}_\beta}$$

and

$$\|v(0,\cdot)\|_{H^{-\frac{1}{3}}(0,\beta)} \leq \|v(0,\cdot)\|_{L^2(0,\beta)} \leq \beta^{\frac{2}{3}} \|v(0,\cdot)\|_{L^6(0,\beta)} \leq C\beta^{\frac{2}{3}} \|v(0,\cdot)\|_{H^{\frac{1}{3}}(0,\beta)} \leq C\beta^{\frac{2}{3}} \|v\|_{\mathcal{X}_\beta},$$

the system (2.8) defines a map  $\Gamma$  from the space  $\mathcal{X}_{\beta}$  to  $\mathcal{X}_{\beta}$  for any  $0 < \beta \leq \max\{1, T\}$  as follows

$$w = \Gamma(v)$$
 for any  $v \in \mathcal{X}_{\beta}$ 

where  $w \in \mathcal{X}_{\beta}$  is the corresponding solution of (2.8) and

$$\|\Gamma(v)\|_{\mathcal{X}_{\beta}} \le C_1 \|\varphi_0\|_{L^2(0,L)} + C_2 \beta^{\frac{1}{2}} \|v\|_{\mathcal{X}_{\beta}}$$

where  $C_1$  and  $C_2$  are two constants depending only on T. Choose r > 0 and  $0 < \beta \le \max\{1, T\}$  such that

$$r = 2C_1 \|\varphi_0\|_{L^2(0,L)}, \quad 2C_2 \beta^{\frac{1}{2}} \le \frac{1}{2}.$$

Then, for any

$$v \in B_{\beta,r} = \{ v \in \mathcal{X}_{\beta}; \quad ||v||_{\mathcal{X}_{\beta}} \le r \},$$

we have

$$\|\Gamma(v)\|_{\mathcal{X}_{\beta}} \leq r.$$

Moreover, for any  $v_1, v_2 \in B_{\beta,r}$ , we get

$$\|\Gamma(v_1) - \Gamma(v_2)\|_{\mathcal{X}_{\beta}} \le 2C_2\beta^{\frac{1}{2}} \|v_1 - v_2\|_{\mathcal{X}_{\beta}} \le \frac{1}{2} \|v_1 - v_2\|_{\mathcal{X}_{\beta}}.$$

Therefore the map  $\Gamma$  is a contraction mapping on  $B_{\beta,r}$ . Its fixed point  $w = \Gamma(w) \in \mathcal{X}_{\beta}$  is the desired solution for  $t \in (0,\beta)$ . As the chosen  $\beta$  is independent of  $\varphi_0$ , the standard continuation extension argument yields that the solution w belongs to  $\mathcal{X}_T$ . The proof is complete.

Finally we conclude this section with an elementary estimate for solutions of system (2.6).

**Proposition 2.9.** Any solution  $\varphi$  of the adjoint problem (2.6) with initial data  $\varphi_0 \in L^2(0,L)$  satisfies

$$\|\varphi_0\|_{L^2(0,L)}^2 \le \frac{1}{T} \|\varphi\|_{L^2((0,T)\times(0,L))}^2 + \|\varphi_x(0,\cdot)\|_{L^2(0,T)}^2 + \|\varphi(0,\cdot)\|_{L^2(0,T)}^2.$$

**Proof.** Multiplying both sides of the equation in (2.6) by  $(T-t)\varphi$  and integrating by parts over  $(0,L)\times(0,T)$ , we get

$$\int_{0}^{L} (T-t)\varphi^{2} \Big|_{0}^{T} dx + \int_{0}^{T} (T-t) \left(\varphi^{2}(0,t) + \varphi_{x}^{2}(0,t)\right) dt + \int_{0}^{T} \int_{0}^{L} \varphi^{2} dx dt = 0.$$

Consequently,

$$\int_{0}^{L} \varphi_{0}^{2} dx \leq \frac{1}{T} \int_{0}^{L} \int_{0}^{T} \varphi^{2} dt dx + \int_{0}^{T} \varphi^{2}(0, t) dt + \int_{0}^{T} \varphi_{x}^{2}(0, t) dt.$$

Equivalently, the following estimate holds for solutions  $\psi$  of the system (2.5):

$$\|\psi_T\|_{L^2(0,L)}^2 \le \frac{1}{T} \|\psi\|_{L^2((0,T)\times(0,L))}^2 + \|\psi_x(L,\cdot)\|_{L^2(0,T)}^2 + \|\psi(L,\cdot)\|_{L^2(0,T)}^2. \tag{2.9}$$

As a comparison, it is worth pointing out that for the adjoint system of (2.3), which is given by

$$\begin{cases} \nu_t + \nu_x + \nu_{xxx} = 0, & \nu(x, T) = \nu_T(x), & (x, t) \in (0, L) \times (0, T), \\ \nu(L, t) = 0, & \nu_x(0, t) = 0, & \nu(0, t) = 0, \end{cases}$$
 (2.10)

the following inequality holds

$$\|\nu_T\|_{L^2(0,L)}^2 \le \frac{1}{T} \|\nu\|_{L^2((0,T)\times(0,L))}^2 + \|\nu_x(L,\cdot)\|_{L^2(0,T)}^2. \tag{2.11}$$

The extra term  $\|\psi(L,\cdot)\|_{L^2(0,T)}^2$  in (2.9) brings new challenges in establishing the observability of the adjoint system (2.5).

## 3 Linear control systems

Consideration is first given to boundary controllability of the linear system

$$\begin{cases} y_t + y_x + y_{xxx} = 0, & (x,t) \in (0,L) \times (0,T), \\ y(0,t) = 0, \ y_x(L,t) = h_2(t), \ y_{xx}(L,t) = 0, \quad t \in (0,T), \end{cases}$$
(3.1)

which employs only one control input  $h_2 \in L^2(0,T)$ .

**Proposition 3.1.** Let  $L \notin \mathbb{F}$  (see (1.18)) and T > 0 be given. There exists a bounded linear operator

$$\Psi:L^2(0,L)\times L^2(0,L)\to L^2(0,T)$$

such that any  $y_0$ ,  $y_T \in L^2(0,L)$ , if one chooses  $h_2 = \Psi(y_0,y_T)$ , then system (3.1) admits a solution  $y \in X_T$  satisfying

$$y|_{t=0} = y_0, \qquad y|_{t=T} = y_T.$$

As it is well-known, the exact controllability of the system (3.1) is related to the observability of its adjoint system

$$\begin{cases} \psi_t + \psi_x + \psi_{xxx} = 0, & (x,t) \in (0,L) \times (0,T), \\ \psi(0,t) = 0, & \psi_x(0,t) = 0, & \psi(L,t) + \psi_{xx}(L,t) = 0, & t \in (0,T), \\ \psi(x,T) = \psi_T(x), & x \in (0,L). \end{cases}$$
(3.2)

**Lemma 3.2.** For all T > 0 and all  $L \notin \mathbb{F}$  there exists C = C(L,T) > 0 such that for any  $\psi_T \in L^2(0,L)$ , the solution  $\psi$  of (3.2) satisfies

$$\|\psi_T\|_{L^2(0,L)} \le C\|\psi_x(L,t)\|_{L^2(0,T)}. (3.3)$$

**Proof.** Proceeding as in [20], if (3.3) is false, then there exists a sequence  $\{\psi_T^n\}_{n\in\mathbb{N}}\in L^2(0,L)$  with  $\|\psi_T^n\|_{L^2(0,L)}=1$  such that the corresponding solutions of (3.2) satisfy

$$1 = \|\psi_T^n\|_{L^2(0,L)} > n\|\psi_T^n(L,\cdot)\|_{L^2(0,T)}.$$

Thus  $\|\psi_x^n(L,\cdot)\|_{L^2(0,T)} \to 0$  as  $n \to \infty$ . By Proposition 2.5 and Theorem 2.6, the sequences  $\{\psi^n\}_{n\in\mathbb{N}}$  and  $\{\psi^n(L,t)\}_{n\in\mathbb{N}}$  are bounded in  $L^2(0,T;H^1(0,L))$  and  $H^{\frac{1}{3}}(0,T)$ , respectively. In addition, according to Proposition 2.9

$$\|\psi_T^n\|_{L^2(0,L)}^2 \leq \frac{1}{T} \|\psi^n\|_{L^2(0,T;L^2(0,L))}^2 + \|\psi_x^n(L,\cdot)\|_{L^2(0,T)}^2 + \|\psi^n(L,\cdot)\|_{L^2(0,T)}^2.$$
 (3.4)

Since,  $\psi_t^n = -(\psi_x^n + \psi_{xxx}^n)$  is bounded in  $L^2(0,T;H^{-2}(0,L))$  and by the embedding

$$H^1(0,L) \hookrightarrow L^2(0,L) \hookrightarrow H^{-2}(0,L)$$

the sequence  $\{\psi^n\}_{n\in\mathbb{N}}$  is relatively compact in  $L^2(0,T;L^2(0,L))$  (see [24]). Furthermore, the second term on the right in (3.4) converges to zero in  $L^2(0,T)$ , and by the compact embedding

$$H^{\frac{1}{3}}(0,T) \hookrightarrow L^2(0,T)$$

the sequence  $\{\psi^n(L,\cdot)\}_{n\in\mathbb{N}}$  has a convergent subsequence on  $L^2(0,T)$ . Therefore  $\{\psi^n_T\}_{n\in\mathbb{N}}$  is a  $L^2(0,L)$ -Cauchy sequence. Let us denote  $\psi_T=\lim_{n\to\infty}\psi^n_T$  and  $\psi$  be the corresponding solution of (3.2). Since  $\psi^n_x(L,t)\to\psi_x(L,t)$  as  $n\to\infty$  in  $L^2(0,T)$  and  $\|\psi^n_T\|_{L^2(0,L)}=1$  for any n, we have  $\|\psi_T\|_{L^2(0,L)}=1$  and  $\psi_x(L,t)=0$ . By the following Lemma 3.3, one can conclude that  $\psi\equiv 0$ , therefore  $\psi_T(x)\equiv 0$  which contradicts the fact that  $\|\psi_T\|_{L^2(0,L)}=1$ .

**Lemma 3.3.** For given T > 0, let us define

 $N_T = \{ \psi_T \in L^2(0, L) : \psi \in X_T \text{ is the mild solution of (3.2) satisfying } \psi_x(L, \cdot) = 0 \text{ in } L^2(0, T) \}.$ Then,  $N_T = \{0\}$  if and only if  $L \notin \mathbb{F}$ . **Proof.** The proof uses the same arguments as that given in [20] and will be presented in the Appendix for the convenience of the interested readers.

Now we turn to prove Proposition 3.1.

**Proof of Proposition 3.1**. Without loss of generality, we assume that  $y_0 = 0$ . Let  $\psi$  be a solution of the system (3.2) and multiply both sides of the equation in (3.1) by  $\psi$  and integrate over the domain  $(0, L) \times (0, T)$ . Integration by parts lead to

$$\int_0^L y(x,T)\psi_t(x)dx = \int_0^T h_2(t)\psi_x(L,t)dt.$$

Let us denote by  $\Upsilon$  the linear and bounded map from  $L^2(0,L) \to L^2(0,L)$  defined by

$$\Upsilon: \psi_T(\cdot) \to y(\cdot, T)$$

with y being the solution of (3.1) when  $h_2(t) = \psi_x(L,t)$  where  $\psi$  is the solution of the system (3.2). According to Lemma 3.2,

$$(\Upsilon(\psi_T), \psi_T)_{L^2(0,L)} = \|\psi_x(L, \cdot)\|_{L^2(0,T)}^2 \ge C^{-2} \|\psi_T\|_{L^2(0,L)}^2. \tag{3.5}$$

Thus  $\Upsilon$  is invertible by Lax-Milgram Theorem. Consequently, for given  $y_T \in L^2(0, L)$ , we can define  $\psi_T = \Upsilon^{-1}y_T$ . We solve system (3.2) and get  $\psi \in X_T$ . Then, we set  $h_2(t) = \psi_x(L, t)$  in system (3.1) and see that the corresponding solution  $y \in X_T$  satisfies

$$y|_{t=0} = 0, \quad y|_{t=T} = y_T.$$

The proof is complete. ■

Next we turn to consider boundary controllability of the linear system

$$\begin{cases} y_t + y_x + y_{xxx} = 0, & (x,t) \in (0,L) \times (0,T), \\ y(0,t) = h_1(t), & y_x(L,t) = h_2(t), & y_{xx}(L,t) = 0, & t \in (0,T), \end{cases}$$
(3.6)

with two control inputs  $h_1 \in H^{\frac{1}{3}}(0,T)$  and  $h_2 \in L^2(0,T)$ .

**Proposition 3.4.** Let T > 0 be given. There exists a bounded linear operator

$$F: L^2(0,L) \times L^2(0,L) \to H^{\frac{1}{3}}(0,T) \times L^2(0,T)$$

such that for any  $y_0, y_T \in L^2(0,L)$ , if one chooses

$$(h_1, h_2) = F(y_0, y_T),$$

then the system (3.6) admits a solution  $y \in X_T$  satisfying

$$y|_{t=0} = y_0, \qquad y|_{t=T} = y_T.$$

As before, we first establish the following observability estimate for the corresponding adjoint system (3.2).

**Lemma 3.5.** Let T > 0 be given. There exists a constant C > 0 such that for any  $\psi_T \in L^2(0, L)$ , the corresponding solution  $\psi$  of (3.2) satisfies

$$\|\psi_T\|_{L^2(0,L)} \le C\left(\int_0^T (|\Delta_t^{-\frac{1}{3}}\psi_{xx}(0,t)|^2 + |\psi_x(L,t)|^2)dt\right)$$
(3.7)

where  $\Delta_t := I - \partial_t^2$ .

**Proof.** If the estimate (3.7) is false, then there exists a sequence  $\{\psi_T^n\}_{n\in\mathbb{N}}\in L^2(0,L)$  with  $\|\psi_T^n\|_{L^2(0,L)}=1$  such that the corresponding solutions  $\psi^n$  of (3.2) satisfies

$$1 = \|\psi_T^n\|_{L^2(0,L)} > n \left( \int_0^T (|\Delta_t^{-\frac{1}{3}} \psi_{xx}^n(0,t)|^2 + |\psi_x^n(L,t)|^2) dt \right)$$

for any n. Thus

$$\|\Delta_t^{-\frac{1}{3}}\psi_{xx}^n(0,\cdot)\|_{L^2(0,T)} \to 0 \quad \text{and} \quad \|\psi_x^n(L,\cdot)\|_{L^2(0,T)} \to 0$$
 (3.8)

when  $n \to \infty$ . Arguing as in the proof of Lemma 3.2 we can conclude that  $\{\psi_T^n\}_{n \in \mathbb{N}}$  is a Cauchy sequence in  $L^2(0,L)$  converging to some  $\psi_T \in L^2(0,L)$ . The corresponding solution  $\psi$  of (3.2) satisfies  $\psi_{xx}(0,t) = 0$  and  $\psi_x(L,t) = 0$ , i.e.,

$$\begin{cases}
\psi_t + \psi_x + \psi_{xxx} = 0, & (x,t) \in (0,L) \times (0,T), \\
\psi(0,t) = 0, & \psi_x(0,t) = 0, & t \in (0,T), \\
\psi(L,t) + \psi_{xx}(L,t) = 0, & \psi_x(L,t) = 0, & t \in (0,T), \\
\psi(x,T) = \psi_T(x), & x \in (0,L),
\end{cases}$$
(3.9)

from which we have  $\psi \equiv 0$  because of the unique continuation property  $(\psi(0,t) = \psi_x(0,t) = \psi_{xx}(0,t) = 0$  for any  $t \in (0,T)$ ). In particular,  $\psi_T \equiv 0$  which contradicts the fact that  $\|\psi_T\|_{L^2(0,L)} = 1$ .

**Proof of Proposition 3.4**: Without loss of generality, we assume that  $y_0 = 0$ . Let  $\psi$  be a solution of the system (3.2) and multiply both sides of the equation in (3.6) by  $\psi$  and integrate over the domain  $(0, L) \times (0, T)$ . Integration by parts leads to

$$\int_0^L y(x,T)\psi_T(x)dx = \int_0^T (h_1(t)\psi_{xx}(0,t) + h_2(t)\psi_x(L,t)) dt.$$

Let us denote by  $\Upsilon$  the linear and bounded map from  $L^2(0,L) \to L^2(0,L)$  defined by

$$\Upsilon: \psi_T(\cdot) \to y(\cdot, T)$$

with y being the solution of (3.6) when

$$h_1(t) = \Delta_t^{-\frac{1}{3}} \psi_{xx}(0, t)$$
 and  $h_2(t) = \psi_x(L, t)$ 

where  $\psi$  is the solution of system (3.2). Thus

$$(\Upsilon(\psi_T), \psi_T)_{L^2(0,L)} = (\Delta_t^{-\frac{1}{3}} \psi_{xx}(0,\cdot), \Delta_t^{-\frac{1}{3}} \psi_{xx}(0,\cdot)_{L^2(0,T)} + \|\psi_x(L,\cdot)\|_{L^2(0,L)}^2 \ge C^{-2} \|\psi_T\|_{L^2(0,L)}^2.$$

The proof is then completed by using the Lax-Milgram Theorem. ■

We now consider boundary controllability of the linear system

$$\begin{cases} y_t + y_x + y_{xxx} = 0, & (x,t) \in (0,L) \times (0,T), \\ y(0,t) = 0, & y_x(L,t) = h_2(t), & y_{xx}(L,t) = h_3(t), & t \in (0,T), \end{cases}$$
(3.10)

with two control inputs  $h_2 \in L^2(0,T)$  and  $h_3 \in H^{-\frac{1}{3}}(0,T)$ .

**Proposition 3.6.** Let T > 0 be given. There exists a bounded linear operator

$$F_1: L^2(0,L) \times L^2(0,L) \to L^2(0,T) \times H^{-\frac{1}{3}}(0,T)$$

such that for any  $y_0, y_T \in L^2(0,L)$ , if one chooses

$$(h_2, h_3) = F_1(y_0, y_T),$$

then the system (3.6) admits a solution  $y \in X_T$  satisfying

$$y|_{t=0} = y_0, \qquad y|_{t=T} = y_T.$$

As before, Proposition 3.6 follows from the following observability estimates for the corresponding adjoint system (3.2).

**Lemma 3.7.** Let T > 0 be given. There exists a constant C > 0 such that for any  $\psi_T \in L^2(0, L)$ , the corresponding solution  $\psi$  of (3.2) satisfies

$$\|\psi_T\|_{L^2(0,L)} \le C\left(\int_0^T (|\Delta_t^{\frac{1}{3}}\psi(L,t)|^2 + |\psi_x(L,t)|^2)dt.\right)$$
(3.11)

**Proof.** The proof is similar to that of Lemma 3.5 and is therefore omitted. ■

Finally we turn to consider the linear system associated to (1.2) using only  $g_1 \in H^{\frac{1}{3}}(0,T)$  and  $g_3 \in L^2(0,T)$  as control inputs, i.e.

$$\begin{cases} u_t + u_x + u_{xxx} = 0, & (x,t) \in (0,L) \times (0,T), \\ u(0,t) = g_1(t), & u(L,t) = 0, & u_x(L,t) = g_3(t), & t \in (0,T). \end{cases}$$
(3.12)

The critical length phenomenon will not occur and system (3.12) is exactly controllable for any L > 0 as stated in the following result.

**Proposition 3.8.** Let T > 0 be given. There exists a bounded linear operator

$$F_2: L^2(0,L) \times L^2(0,L) \to H^{\frac{1}{3}}(0,T) \times L^2(0,T)$$

such that for any  $u_0$ ,  $u_T \in L^2(0,L)$ , if one chooses

$$(g_1, g_3) = \digamma_2(u_0, u_T),$$

then the system (3.12) admits a solution  $y \in X_T$  satisfying

$$u|_{t=0} = u_0, \qquad u|_{t=T} = u_T.$$

**Proof.** The proof is similar to that of Proposition 3.6 and is therefore skipped.

## 4 Nonlinear control systems

In this section we first consider the nonlinear system

$$\begin{cases} y_x + y_{xxx} + y_x + yy_x = 0, & x \in (0, L), \ t \in (0, T), \\ y(0, t) = 0, \ y_x(L, t) = h_2(t), \ y_{xx}(L, t) = 0, \quad t \in (0, T), \\ y(x, 0) = y_0(x), \quad x \in (0, L), \end{cases}$$

$$(4.1)$$

and present the proof of Theorem 1.1.

**Proof of Theorem 1.1.** Rewrite the system (4.1) in its integral form

$$y(t) = W_0(t)y_0 + W_{bdr}(t)h_2 - \int_0^t W_0(t - \tau)(yy_x)(\tau)d\tau.$$
 (4.2)

Here we have written  $W_{bdr}(t)(0, h_2, 0)$  as  $W_{bdr}(t)h_2$  for simplicity. For any  $v \in X_T$ , let us set

$$\nu(T,v) := \int_0^T W_0(T-\tau)(vv_x)(\tau)d\tau.$$

For any  $y_0, y_T \in L^2(0, L)$ , we use Proposition 3.1 to define

$$h_2 = \Psi(y_0, y_T + \nu(T, v)).$$

Then

$$v(t) = W_0(t)y_0 + W_{bdr}\Psi(y_0, y_T + \nu(T, v)) - \int_0^t W_0(t - \tau)(vv_x)(\tau)d\tau$$

satisfies

$$v|_{t=0} = y_0, \qquad v|_{t=T} = y_T + \nu(T, v) - \nu(T, v) = y_T.$$

This leads us to consider the map

$$\Gamma(v) = W_0(t)y_0 + W_{bdr}\Psi(y_0, y_T + \nu(T, v)) - \int_0^t W_0(t - \tau)(vv_x)(\tau, x)d\tau.$$

If we can show that the map  $\Gamma$  is a contraction in an appropriate metric space, then its fixed point v is a solution of (4.1) with  $h_2 = \Psi(y_0, y_T + \nu(T; v))$  which satisfies

$$v|_{t=0} = y_0, \qquad v|_{t=T} = y_T.$$

Next we show that this is indeed the case; the map  $\Gamma$  is a contraction map in the ball

$$B_r = \{ z \in X_T; \ \|z\|_{X_T} \le r \}$$

for an appropriately chosen r. According to Proposition 2.1, there exists a constant  $C_1 > 0$  such that

$$\|\Gamma(v)\|_{X_T} \le C_1 \left( \|y_0\|_{L^2(0,L)} + \|\Psi(y_0, y_T + \nu(T,v))\|_{L^2(0,L)} + \int_0^T \|vv_x\|_{L^2(0,L)}(t)dt \right).$$

Since

$$\|\Psi(y_0, y_T + \nu(T, v))\|_{L^2(0, L)} \le C_2 \left( \|y_0\|_{L^2(0, L)} + \|y_T\|_{L^2(0, L)} + \|\nu(T, v)\|_{L^2(0, L)} \right),$$

$$\|\nu(T,v)\|_{L^2(0,L)} \le \int_0^T \|W_0(T-\tau)vv_x\|_{L^2(0,L)} d\tau \le \int_0^T \|vv_x\|_{L^2(0,L)}(t) dt,$$

and the bilinear estimate

$$\int_{0}^{T} \|vv_{x}\|_{L^{2}(0,L)}(t)dt \leq C_{3}\|v\|_{X_{T}}^{2}$$

we arrive at

$$\|\Gamma(v)\|_{X_T} \le C_3(\|y_0\|_{L^2(0,L)} + \|y_T\|_{L^2(0,L)}) + C_4\|v\|_{X_T}^2$$

for any  $v \in X_T$  where  $C_3$  and  $C_4$  are constants depending only on T. By choosing  $r, \delta$  such that

$$r = 2C_3\delta, \qquad 4C_3C_4\delta < \frac{1}{2}$$
 (4.3)

we get

$$\|\Gamma(v)\|_{X_T} \le C_3\delta + 4C_4C_3\delta C_3\delta \le 2C_3\delta \le r$$

for any  $v \in B_r$ . In addition, for  $v_1, v_2 \in B_r$ ,

$$\Gamma(v_1) - \Gamma(v_2) = W_{bdr} \Psi \left( 0, \nu(T, v_1) - \nu(T, v_2) \right) + \frac{1}{2} \int_0^t W_0(t - \tau) \left[ (v_1 + v_2)(v_1 - v_2)_x \right] (\tau) d\tau$$

and

$$\begin{split} \|\Gamma(v_1) - \Gamma(v_2)\|_{X_T} & \leq & \left\{ C_4(\|v_1\|_{X_T} + \|v_2\|_{X_T}) + C_4(\|v_1\|_{X_T} + \|v_2\|_{X_T}) \right\} \|v_1 - v_2\|_{X_T} \\ & \leq & 8C_3C_4\delta\|v_1 - v_2\|_{X_T} \\ & \leq & \alpha\|v_1 - v_2\|_{X_T} \end{split}$$

with  $\alpha = 8C_3C_4\delta < 1$ . The proof is completed.

As Theorem 1.4 and Theorem 1.5 can be proved using the same arguments as those in the proof of Theorem 1.1, their proofs will be skipped.

Now we turn to consider the system

$$\begin{cases} y_x + y_{xxx} + y_x + yy_x = 0, & x \in (0, L), \ t \in (0, T), \\ y(0, t) = 0, \ y_x(L, t) = 0, \ y_{xx}(L, t) = h_3(t), & t \in (0, T) \end{cases}$$

$$(4.4)$$

to prove Theorem 1.3.

**Proof of Theorem 1.3.** We first consider the system

$$\begin{cases}
 u_t + u_x + u_{xxx} + uu_x = 0, & (x,t) \in (0,L) \times (0,T), \\
 u(0,t) = 0, & u(L,t) = g_2(t), & u_x(L,t) = 0, & t \in (0,T),
\end{cases}$$
(4.5)

and show that the following controllability result holds, which is an improvement of Theorem C due to Glass and Guerrero [10].

**Proposition 4.1.** Let T > 0 and  $L \notin \mathcal{N}$  be given. There exists r > 0 such that for any  $u_0, u_T \in L^2(0, L)$  with

$$||u_0||_{L^2(0,L)} + ||u_T||_{L^2(0,L)} \le r,$$

there exists  $g_2 \in H^{\frac{1}{3}}(0,T)$  such that the system (4.5) admits a solution

$$u \in C([0,T]; L^2(0,L)) \cap L^2(0,T; H^1(0,L))$$

satisfying

$$u|_{t=0} = u_0, \qquad u|_{t=T} = u_T.$$

If Proposition 4.1 holds, for given  $y_0, y_T \in L^2(0, L)$ , set

$$u_0 = y_0, \quad u_T = y_T.$$

Then by Proposition 4.1, there exists  $g_2 \in H^{\frac{1}{3}}(0,T)$  such that (4.5) admits a unique solution  $u \in X_T$  satisfying

$$u|_{t=0} = y_0, \qquad u|_{t=T} = y_T.$$

Thus y(x,t) := u(x,t) will be a desired solution of (4.4) with  $h_3(t) = u_{xx}(L,t)$  satisfying

$$y|_{t=0} = y_0, \quad y|_{t=T} = y_T.$$

As  $u \in X_T$  solves

$$\left\{ \begin{array}{l} u_t + u_x + u_{xxx} = f, \quad u(x,0) = y_0, \quad (x,y) \in (0,L) \times (0,T), \\ u(0,t) = 0, \quad u(L,t) = g_2(t), \quad u_x(L,t) = 0, \quad t \in (0,T), \end{array} \right.$$

with  $g_2 \in H^{\frac{1}{3}}(0,T)$  and  $f = -uu_x \in L^1(0,T;L^2(0,L))$ , it follows from Proposition 2.3 that  $u_{xx}(L,\cdot) \in H^{-\frac{1}{3}}(0,T)$ . Thus, it suffices to prove Proposition 4.1 to complete the proof of Theorem 1.3.

To this end, note that, according to Theorem C,  $g_2 \in H^{\frac{1}{5}-\epsilon}(0,T)$ . We just need to prove that this  $g_2$  given by Theorem C belongs, in fact, to the space  $H^{\frac{1}{3}}(0,T)$ . Indeed, the solution  $u \in X_T$  given in Theorem C can be written as

$$u = \kappa + \mu$$

where  $\kappa$  solves

$$\begin{cases} \kappa_t + \kappa_{xxx} = f, & \kappa(x,0) = u_0, \ (x,t) \in (0,L) \times (0,T), \\ \kappa(0,t) = \kappa(L,t) = \kappa_x(L,t) = 0, & t \in (0,T), \end{cases}$$

with  $f = -u_x - uu_x$ , and  $\mu$  solves

$$\begin{cases}
\mu_t + \mu_{xxx} = 0, & \mu(x,0) = 0, (x,t) \in (0,L) \times (0,T), \\
\mu(0,t) = 0, & \mu(L,t) = g_2(t), & \mu_x(L,t) = 0, t \in (0,T).
\end{cases}$$
(4.6)

As  $u_0 \in L^2(0, L)$  and  $f = -u_x - uu_x \in L^1(0, T; L^2(0, L))$  (because  $u \in X_T$ ), we have  $\kappa \in X_T$  by Proposition 2.3. In addition, the following lemma (whose proof will be presented in the Appendix) holds for system (4.6).

**Lemma 4.2.** The solution  $\mu$  of (4.6) belongs to  $X_T$  if and only if  $g_2$  belongs to  $H^{\frac{1}{3}}(0,T)$ .

Using this Lemma, we get  $u \in X_T$  if and only if  $g_2 \in H^{\frac{1}{3}}(0,T)$ . The proof of Theorem 1.3 is complete.

Finally we consider the system

$$\begin{cases} y_x + y_{xxx} + y_x + yy_x = 0, & x \in (0, L), \ t \in (0, T), \\ y(0, t) = h_1(t), \ y_x(L, t) = h_2(t), \ y_{xx}(L, t) = h_3(t), & t \in (0, T), \\ y(x, 0) = y_0(x), & x \in (0, L), \end{cases}$$

$$(4.7)$$

and present the proof of Theorem 1.7.

**Proof of Theorem 1.7.** Consider first the following initial value control problem for the KdV equation posed on the whole line  $\mathbb{R}$ 

$$\begin{cases} z_t + z_x + zz_x + z_{xxx} = 0, & x, \ t \in \mathbb{R}, \\ z(x,0) = h(x) \end{cases}$$
 (4.8)

where the initial value h(x) is considered as a control input. The following result is due to Zhang [26].

**Theorem G.** Let  $s \geq 0$  and T > 0 be given and suppose  $w \in C(\mathbb{R}; H^{\infty}(\mathbb{R}))$  is a given solution of

$$w_t + w_x + ww_x + w_{xxx} = 0, \quad x, t \in \mathbb{R}.$$

There exists  $\delta > 0$  such that for any  $y_0, y_T \in H^s(0, L)$  with

$$||y_0 - w(\cdot, 0)||_{H^s(0,L)} \le \delta, \qquad ||y_T - w(\cdot, T)||_{H^s(0,L)} \le \delta,$$
 (4.9)

one can find a control input  $h \in H^s(\mathbb{R})$  which is an external modification of  $y_0$  such that (4.8) admits a solution  $z \in C(\mathbb{R}; H^s(\mathbb{R}))$  satisfying

$$z(x,0) = y_0(x),$$
  $z(x,T) = y_T(x)$  for any  $x \in (0,L)$ .

Let u be as in Theorem 1.7. Applying Theorem G with w = u and s = 0, we get the existence of  $\delta > 0$  such that for any  $y_0, y_T \in L^2(0, L)$  with

$$||y_0 - u(\cdot, 0)||_{L^2(0, L)} \le \delta, \qquad ||y_T - u(\cdot, T)||_{L^2(0, L)} \le \delta,$$
 (4.10)

there exists  $h \in L^2(\mathbb{R})$  and the corresponding  $z \in C(\mathbb{R}; L^2(\mathbb{R})) \cap L^2(\mathbb{R}; H^1_{loc}(\mathbb{R}))$  solution of (4.8). Let y be the restriction of z to the domain  $(0, L) \times (0, T)$ , and

$$h_1(t) = z(0,t), \quad h_2(t) = z_x(L,t), \quad h_3(t) = z_{xx}(L,t)$$

for 0 < t < T. Then, according to [28], we have that  $y \in X_T$  solves (4.7) with  $h_1 \in H^{\frac{1}{3}}(0,T)$ ,  $h_2 \in L^2(0,T)$  and  $h_3 \in H^{-\frac{1}{3}}(0,T)$ . Moreover,

$$y|_{t=0} = y_0, \qquad y|_{t=T} = y_T.$$

The proof is complete.  $\blacksquare$ 

## 5 Conclusion Remarks

The focus of our discussion has been on the boundary controllability of two classes of boundary control systems described by the KdV equation posed on a finite domain (0, L), namely,

$$\begin{cases}
 u_t + u_x + uu_x + u_{xxx} = 0, & (x,t) \in (0,L) \times (0,T), \\
 u(0,t) = g_1(t), & u(L,t) = g_2(t), & u_x(L,t) = g_3(t), & t \in (0,T),
\end{cases}$$
(5.1)

and

$$\begin{cases} y_t + y_x + yy_x + y_{xxx} = 0, & (x,t) \in (0,L) \times (0,T), \\ y(0,t) = h_1(t), & y_x(L,t) = h_2(t), & y_{xx}(L,t) = h_3(t) & t \in (0,T). \end{cases}$$
(5.2)

The linear systems associated to these equations are obtained by dropping the nonlinear term  $uu_x$  and  $yy_x$ , respectively. The system (5.1) has been intensively studied and various controllability results have been established in the past. However, there have been few results for the second system (5.2) because of some difficulties to apply directly the methods that work effectively for the system (5.1). In this paper, aided by the newly established hidden regularities of solutions of the KdV equation, we have succeeded in overcoming those difficulties and established various boundary controllability results for the system (5.2) similar to those known for the system (5.1) in the literature. Furthermore, with the new tool in hand, we have also be able to improve some known controllability results for the system (5.1). Our results can be summarized as below.

- (i) The linear system associated to (5.2) is exactly controllable with two or three boundary controls in action. In any of those cases, the nonlinear system (5.2) is locally exactly controllable.
- (ii) With only a single control  $h_2$  in action ( $h_1 = h_3 = 0$ ), the linear system associated to (5.2) is exactly controllable if and only if L does not belong to

$$\mathbb{F} = \left\{ L \in \mathbb{R}^+ : L^2 = -(a^2 + ab + b^2) \text{ with } a, b \in \mathbb{C} \text{ satisfying} \quad \frac{e^a}{a^2} = \frac{e^b}{b^2} = \frac{e^{-(a+b)}}{(a+b)^2} \right\}.$$

The nonlinear system (5.2) is also locally exactly controllable if  $L \notin \mathbb{F}$ .

(iii) When only control input  $h_3$  is employed ( $h_1 = h_2 = 0$ ), the linear system associated to (5.2) is exactly controllable if and only if L does not belong to

$$\mathcal{N} = \Big\{ L \in \mathbb{R}^+ : L^2 = -(a^2 + ab + b^2) \text{ with } a,b \in \mathbb{C} \text{ satisfying} \quad ae^a = be^b = -(a+b)e^{-(a+b)} \Big\}.$$

Moreover, if  $L \notin \mathcal{N}$ , then the nonlinear system (5.2) is locally exactly controllable.

- (iv) The linear system associated to (5.1) is exactly controllable with control inputs  $g_1$  and  $g_2$  in action (put  $g_3 = 0$ ). In this case, the nonlinear system (5.1) is locally exactly controllable.
- (v) We have improved the regularity of the control  $g_2$  in (5.1). In previous work [10], the control  $g_2$  is known to belong to the space  $H^{\frac{1}{6}-\epsilon}(0,T)$  for any  $\epsilon > 0$ . In this paper we are able to prove that the control input  $g_2$  belongs in fact to the space  $H^{\frac{1}{3}}(0,T)$ .

While some significant progresses have been made in the study of boundary controllability of the KdV equation on a bounded domain, there are still a lot of interesting questions left open for further investigations. One of them is the so-called critical length problem. As it is well known now, the linear systems associated to (5.1) and (5.2) are not always exactly controllable if only a single control input is allowed to act on the right end of the spatial domain (0, L). In general, if the associated linear system is not exactly controllable, one would intend to believe the nonlinear system is also not exactly controllable. However, for the system (5.1) with only control input  $g_3$  in action, though its associated linear system is not exactly controllable when  $L \in S$  (see (1.5) for the definition of S), the nonlinear system (5.1) has been shown by Coron and Crepeau [8], Cerpa [3], and Cerpa and Creapeau [4] to be locally (large time) exactly controllable. The questions still remain open for other critical length problems.

## Open Problem 5.1. (Critical length problems)

(a) Is the nonlinear system (5.2) with only control input  $h_2$  in action exactly controllable when the length L of the spatial domain (0, L) belongs to the set  $\mathbb{F}$ ?

- (b) Is the nonlinear system (5.2) with only control input  $h_3$  in action exactly controllable when the length L of the spatial domain (0,L) belongs to the set  $\mathcal{N}$ ?
- (c) Is the nonlinear system (5.1) with only control input  $g_3$  in action exactly controllable when the length L of the spatial domain (0,L) belongs to the set  $\mathcal{N}$ ?

Most controllability results that have been established so far for both systems (5.1) and (5.2) are local: one can only guide a small amplitude initial state to a small amplitude terminal state by choosing appropriate boundary control inputs. The following question arises naturally.

**Open Problem 5.2.** (Global controllability problem) Are the nonlinear systems (5.1) and (5.2) globally exactly boundary controllable?

The following global interior stabilization result for the KdV equation on a finite interval

$$\begin{cases}
 u_t + u_x + uu_x + u_{xxx} + a(x)u = 0, & u(x,0) = u_0(x), & x \in (0,L), t > 0, \\
 u(0,t) = 0, & u(L,t) = 0, & u_x(L,t) = 0, & t > 0
\end{cases}$$
(5.3)

is well-known in the literature (see [17, 18, 22]).

**Theorem G.** Assume the function  $a \in L^{\infty}(0,L)$  with  $a(x) \geq 0$  and such that the support of a is a nonempty open subset of (0,L). There exists  $\gamma > 0$  such that for any  $u_0 \in L^2(0,L)$ , the corresponding solution u of (5.3) belongs to the space  $C([0,\infty);L^2(0,L))$  and, moreover,

$$||u(\cdot,t)||_{L^2(0,L)} \le \alpha(||u_0||_{L^2(0,L)})e^{-\gamma t} \quad \forall \ t \ge 0$$

where  $\alpha: \mathbb{R}^+ \to \mathbb{R}^+$  is a nondecreasing continuous function.

Combining Theorem G, Theorems 1.3,1.4, 1.5 we have the following partial answers to the Open Problem 5.2 for the nonlinear system (5.2).

**Theorem 5.1.** There exists  $\delta > 0$ . For any N > 0, one can find a T > 0 depending only on N and  $\delta$  such that for any  $\phi$ ,  $\psi \in L^2(0,L)$  with

$$\|\phi\|_{2(0,L)} \le N, \qquad \|\psi\|_{2(0,L)} \le \delta,$$

one can find either

$$h_1 \in H^{\frac{1}{3}}(0,T), \quad h_2 \in L^2(0,T), \quad h_3 = 0,$$

or

$$h_1 \in H^{\frac{1}{3}}(0,T), \quad h_2 = 0, \quad h_3 \in H^{-\frac{1}{3}}(0,T),$$

or

$$h_1 = 0, \quad h_2 \in L^2(0,T), \quad h_3 \in H^{-\frac{1}{3}}(0,T),$$

such that the nonlinear system (5.2) admits a solution  $u \in C([0,T];L^2(0,L))$  satisfying

$$u|_{t=0} = \phi, \qquad u|_{t=T} = \psi.$$

*Proof.* Note first that for any  $\phi \in L^2(0,L)$ , the system (5.3) admits a unique solution  $u \in C([0,\infty); L^2(0,L))$  which also possesses the hidden regularity

$$\partial_x^k u(x,t) \in L_x^{\infty}(0,L; H^{\frac{1-k}{3}}(0,T), \quad k = 0,1,2.$$

Then Theorem 5.1 follows from Theorem G and Theorems 1.3, 1.4, 1.5 using the same argument as that used in the proof of Theorem 3.22 in [23].  $\blacksquare$ 

**Remark 5.2.** Theorem 5.1 only provides a partial answer to Problem 5.2 since the amplitude of the terminal state is still required to be small. Question remains:

Can small amplitude restriction on the terminal state be removed?

If one is allowed to use all three boundary control inputs, then the small amplitude restriction can be removed.

**Theorem 5.3.** Let N > 0 be given. There exists a T > 0 such that for any  $\phi$ ,  $\psi \in L^2(0,L)$  with

$$\|\phi\|_{L^2(0,L)} \le N, \qquad \|\psi\|_{L^2(0,L)} \le N,$$

the nonlinear equation

$$u_t + u_x + uu_x + u_{xxx} = 0, \quad x \in (0, L) \times (0, T)$$

admits a solution  $u \in C([0,T]; L^2(0,L))$  satisfying

$$u|_{t=0} = \phi, \qquad u|_{t=T} = \psi.$$

*Proof.* For given  $\phi$ ,  $\psi \in L^2(0,L)$ , let  $\tilde{\phi}$  and  $\tilde{\psi}$  be their extension from (0,L) to (0,2L) such that

$$\tilde{\phi} \in L^2(0, 2L), \qquad \tilde{\psi} \in L^2(0, 2L), \qquad \int_0^{2L} \tilde{\phi}(x) dx = \int_0^{2L} \tilde{\psi}(x) dx$$

and consider the following internal control problem of the KdV equation posed on the interval (0, 2L) with periodic boundary condition

$$\begin{cases} v_t + v_x + vv_x + v_{xxx} + a(x)v = 0, & v(x,0) = \phi(x), & x \in (0,2L), \\ v(0,t) = v(2L,t), & v_x(0,t) = v_x(2L,t), & v_{xx}(0,t) = v_{xx}(2L,t) \end{cases}$$

where  $a \in L^{\infty}(0, 2L)$  and support of  $a \in (L, 2L)$ . The proof is completed by invoking Theorem 1.1 in [15].

Consequently, if chooses

$$h_1(t) = u(0,t),$$
  $h_2(t) = u_x(L,t),$   $h_3(t) = u_{xx}(L,t),$ 

then the system (5.2) will be guided from the given initial state  $\phi$  to the given terminal state  $\psi$ . The only drawback is that we do not know exactly the regularities of the boundary inputs  $h_j$ , j = 1, 2, 3.

In Theorem 5.1 and Theorem 5.3 the time interval (0,T) used to conduct control depends on the size of the initial state and terminal state. The larger of the size of the initial state, the longer the time interval (0,T). Such type of controllability is usually called the large time controllability. As it is well-known, the KdV equation possesses infinite propagation speed. Thus one may wonder the following.

**Open Problem 5.3.** Can the time interval (0,T) in Theorems 5.1 and 5.3 be chosen arbitrarily small?

## 6 Appendices

### 6.1 Proofs of Proposition 2.8 and Lemma 4.2

**Proof of Proposition 2.8.** The solution of the system

$$\begin{cases} w_t + w_{xxx} = f, & w(x,0) = w_0(x), & (x,t) \in (0,L) \times (0,T), \\ w_{xx}(0,t) = k_1(t), & w(L,t) = k_2(t), & w_x(L,t) = k_3(t), & t \in (0,T) \end{cases}$$

$$(6.1)$$

can be written as

$$w(t) = W_0(t)w_0 + W_{bdr}(t)\vec{k} + \int_0^t W_0(t-\tau)f(\tau)d\tau$$

with  $\vec{k} = (k_1, k_2, k_3)$  where  $W_0(t)$  is the  $C_0$  semigroup in  $L^2(0, L)$  generated by the operator

$$Bf = -f'''$$

with the domain

$$\mathcal{D}(B) = \{ f \in H^3(0, L); \ f''(0) = f(L) = f'(L) = 0 \}.$$

Then,  $u(t) = W_0(t)w_0$  solves

$$\begin{cases} u_t + u_{xxx} = 0, & u(x,0) = w_0(x), & (x,t) \in (0,L) \times (0,T), \\ u_{xx}(0,t) = 0, & u(L,t) = 0, & u_x(L,t) = 0, & t \in (0,T), \end{cases}$$

$$(6.2)$$

 $v(t) = W_{bdr}(t)\vec{k}$  solves

$$\begin{cases} v_t + v_{xxx} = 0, & v(x,0) = 0, \quad (x,t) \in (0,L) \times (0,T), \\ v_{xx}(0,t) = k_1(t), & v(L,t) = k_2(t), & v_x(L,t) = k_3(t), \quad t \in (0,T), \end{cases}$$
 (6.3)

and  $z(t) = \int_0^t W_0(t-\tau)f(\tau)d\tau$  solves

$$\begin{cases}
z_t + z_{xxx} = f, & z(x,0) = 0, & (x,t) \in (0,L) \times (0,T), \\
z_{xx}(0,t) = 0, & z(L,t) = 0, & z_x(L,t) = 0, & t \in (0,T).
\end{cases}$$
(6.4)

As in the proof of Proposition 2.1, it is easy to see that for any  $f \in L^1(0,T;L^2(0,L))$  and  $w_0 \in L^2(0,L)$ , both  $u = W_0(t)w_0$  and  $z = \int_0^t W_0(t-\tau)f(\tau)d\tau$  belong to the space  $X_T$  and, in addition, there exists a constant C > 0 such that

$$||u||_{X_T} + ||z||_{X_T} \le C \left( ||w_0||_{L^2(0,L)} + ||f||_{L^1(0,T;L^2(0,L))} \right).$$

For  $v(t) = W_{bdr}(t)\vec{k}$ , following [2], we first look for an explicit representation formula. Applying the Laplace transform with respect to t in both sides of the equation in (6.3), (i.e.  $\hat{v}(s,x) = \int_0^\infty e^{-st}v(t)dt$ ), we obtain

$$\begin{cases} s\hat{v} + \hat{v}_{xxx} = 0, & x \in (0, L), \ s > 0, \\ \hat{v}_{xx}(0, s) = \hat{k}_1(s), & \hat{v}(L, s) = \hat{k}_2(s), & \hat{v}_x(L, s) = \hat{k}_3(s), \quad s > 0. \end{cases}$$

$$(6.5)$$

Its solution  $\hat{v}(x,s)$  can be written as  $\hat{v}(x,s) = \sum_{j=1}^{3} c_j(s) e^{\lambda_j(s)x}$  where  $\lambda_j$  solves characteristic equation  $s + \lambda^3 = 0$ , i.e.

$$\lambda_1 = i\rho, \quad \lambda_2 = -i\rho\left(\frac{1+i\sqrt{3}}{2}\right), \quad \lambda_3 = -i\rho\left(\frac{1-i\sqrt{3}}{2}\right)$$

with  $s = \rho^3$ . Imposition of the boundary conditions of (6.5) yields that  $c_j = c_j(s)$  for j = 1, 2, 3 solves the system

$$\begin{pmatrix} \lambda_1^2 & \lambda_2^2 & \lambda_3^2 \\ e^{\lambda_1 L} & e^{\lambda_2 L} & e^{\lambda_3 L} \\ \lambda_1 e^{\lambda_1 L} & \lambda_2 e^{\lambda_2 L} & \lambda_3 e^{\lambda_3 L} \end{pmatrix} \begin{pmatrix} c_1 \\ c_2 \\ c_3 \end{pmatrix} = \begin{pmatrix} \hat{k}_1 \\ \hat{k}_2 \\ \hat{k}_3 \end{pmatrix}.$$

By Cramer rule,

$$c_j = \frac{\Delta_j}{\Lambda}$$
, for  $j = 1, 2, 3$ ,

where

$$\Delta = \Delta(s) = \begin{vmatrix} \lambda_1^2 & \lambda_2^2 & \lambda_3^2 \\ e^{\lambda_1 L} & e^{\lambda_2 L} & e^{\lambda_3 L} \\ \lambda_1 e^{\lambda_1 L} & \lambda_2 e^{\lambda_2 L} & \lambda_3 e^{\lambda_3 L} \end{vmatrix}$$

and  $\Delta_j(s)$  is the determinant of the matrices obtained by changing the *j*th-column of  $\Delta$  by the vector  $(\hat{k}_1, \hat{k}_2, \hat{k}_3)^T$  for j = 1, 2, 3. Taking the inverse Laplace transform of  $\hat{v}$  and following the

same arguments as those in [1] lead us to the following representation of the solution v of the system (6.3):

$$v(x,t) = \sum_{m=1}^{3} v_m(x,t)$$

with

$$v_m(x,t) = \sum_{j=1}^{3} v_{j,m}(x,t)$$
 and  $v_{j,m}(x,t) = v_{j,m}^{+}(x,t) + v_{j,m}^{-}(x,t)$ 

where for m, j = 1, 2, 3,

$$v_{j,m}^{+}(x,t) = \frac{1}{2\pi} \int_{0}^{\infty} e^{i\rho^{3}t + \lambda_{j}^{+}(\rho)x} \frac{\Delta_{j,m}^{+}(\rho)}{\Delta^{+}(\rho)} \hat{k}_{m}^{+}(\rho) 3\rho^{2} d\rho,$$
$$v_{j,m}^{-}(x,t) = \overline{v_{j,m}^{+}(x,t)}$$

and

$$\hat{k}_{m}^{+}(\rho) = \hat{k}_{m}(i\rho^{3}), \ \Delta^{+}(\rho) = \Delta(i\rho^{3}), \ \Delta^{+}_{j,m}(\rho) = \Delta_{j,m}(i\rho^{3}), \ \lambda^{+}_{j}(\rho) = \lambda_{j}(i\rho^{3}).$$

**Lemma 6.1.** Let T > 0 be given. There exists a constant C > 0 such that for any  $\vec{k} \in \mathcal{K}_T$ , the system (6.3) admits a unique solution  $v \in X_T$ . Moreover, there exists a constant C > 0 such that

$$||v||_{X_T} + \sum_{j=0}^{2} ||\partial_x^j v||_{L_x^{\infty}(0,L;H^{(1-j)/3}(0,T))} \le C||\vec{k}||_{\mathcal{K}_T}.$$

**Proof.** Note that as stated above, the solution v can be written as

$$v(x,t) = v_1(x,t) + v_2(x,t) + v_3(x,t).$$

Let us prove Lemma 6.1 for  $v_1$ . First of all, by straightforward computation, we can list the asymptotic behavior of the ratios  $\frac{\Delta^+_{j,m}(\rho)}{\Delta^+(\rho)}$  for  $\rho \to +\infty$  as below.

$\frac{\Delta_{1,1}^{+}(\rho)}{\Delta^{+}(\rho)} \sim \rho^{-2} e^{-\frac{\sqrt{3}}{2}\rho L}$	$\frac{\Delta_{2,1}^{+}(\rho)}{\Delta^{+}(\rho)} \sim \rho^{-2} e^{-\sqrt{3}\rho L}$	$\frac{\Delta_{3,1}^{+}(\rho)}{\Delta^{+}(\rho)} \sim \rho^{-2} e^{-\sqrt{3}\rho L}$
$\frac{\Delta_{1,2}^{+}(\rho)}{\Delta^{+}(\rho)} \sim 1$	$\frac{\Delta_{2,2}^{+}(\rho)}{\Delta^{+}(\rho)} \sim e^{-\sqrt{3}\rho L}$	$\frac{\Delta_{3,2}^+( ho)}{\Delta^+( ho)}\sim 1$
$\frac{\Delta_{1,3}^+(\rho)}{\Delta^+(\rho)} \sim \rho^{-1}$	$\frac{\Delta_{2,3}^{+}(\rho)}{\Delta^{+}(\rho)} \sim \rho^{-1} e^{-\frac{\sqrt{3}}{2}\rho L}$	$\frac{\Delta_{3,3}^+(\rho)}{\Delta^+(\rho)} \sim \rho^{-1}$

As

$$v_1(x,t) = \frac{3}{\pi} \sum_{i=1}^{3} \mathcal{R}e \int_0^{\infty} e^{i\rho^3 t} e^{\lambda_j^+(\rho)x} \frac{\Delta_{j,1}^+(\rho)}{\Delta^+(\rho)} \hat{k}_1^+(\rho) \rho^2 d\rho,$$

we have

$$\begin{split} \sup_{0 < t < T} \|v_1(\cdot,t)\|_{L^2(0,L)}^2 & \leq & C \int_0^\infty \rho^{-2} |\hat{k_1}^+(\rho)|^2 \rho^2 d\rho \\ & \leq & C \int_0^\infty \mu^{-2/3} |\hat{k_1}(i\mu)|^2 d\mu \\ & \leq & C \|k_1\|_{H^{-\frac{1}{3}}(\mathbb{R}^+)}^2 \\ & \leq & C \|\vec{k}\|_{\mathcal{K}_T}. \end{split}$$

Furthermore, for  $\ell = -1, 0, 1$ , set  $\mu = \rho^3$ ,  $\theta(\mu) = \mu^{\frac{1}{3}}$ ,

$$\begin{array}{lcl} \partial_{x}^{\ell+1}v_{1}(x,t) & = & \frac{3}{\pi}\sum_{j=1}^{3}\mathcal{R}e \int_{0}^{\infty}(\lambda_{j}^{+}(\rho))^{\ell+1}e^{i\rho^{3}t}e^{\lambda_{j}^{+}(\rho)x}\frac{\Delta_{j,1}^{+}(\rho)}{\Delta^{+}(\rho)}\hat{k}_{1}^{+}(\rho)\rho^{2}d\rho\\ & = & \frac{1}{\pi}\sum_{j=1}^{3}\mathcal{R}e \int_{0}^{\infty}(\lambda_{j}^{+}(\theta(\mu)))^{\ell+1}e^{i\mu t}e^{\lambda_{j}^{+}(\theta(\mu))x}\frac{\Delta_{j,1}^{+}(\theta(\mu))}{\Delta^{+}(\theta(\mu))}\hat{k}_{1}(i\mu)d\mu \end{array}$$

Applying Plancherel Theorem in time t yields that, for any  $x \in (0, L)$ ,

$$\begin{split} \|\partial_x^{\ell+1} v_1(x,\cdot)\|_{H^{-\frac{\ell}{3}}(0,T)}^2 & \leq C \sum_{j=1}^3 \int_0^\infty \mu^{-\frac{2\ell}{3}} \left| (\lambda_j^+(\theta(\mu)))^{\ell+1} e^{\lambda_j^+(\theta(\mu))x} \frac{\Delta_{j,1}^+(\theta(\mu))}{\Delta^+(\theta(\mu))} \hat{k}_1(i\mu) \right|^2 d\mu \\ & \leq C \int_0^\infty \mu^{-\frac{2\ell}{3}} |\hat{k}_1(i\mu)|^2 d\mu \\ & \leq C \|k_1\|_{H^{-\frac{\ell}{3}}(0,T)}^2 \\ & \leq C \|\vec{k}\|_{\mathcal{K}_T}^2. \end{split}$$

for  $\ell = -1, 0, 1$ . Consequently

$$\sup_{0 < x < L} \|\partial_x^{\ell+1} v_1(x, \cdot)\|_{H^{-\frac{\ell}{3}}(0, T)} \le C \|\vec{k}\|_{\mathcal{K}_T}$$

for  $\ell = -1, 0, 1$ . In particular,

$$||v_1||_{L^2(0,T;H^1(0,L))} \le C||\vec{k}||_{\mathcal{K}_T}$$

which ends the proof of Lemma 6.1 for  $v_1$ . The proofs for  $v_2$  and  $v_3$  are similar.

Now we turn to complete the proof of Proposition 2.8. It remains to prove that,

$$\|\partial_x^j u\|_{L_x^{\infty}(0,L;H^{\frac{1-j}{3}}(0,T))} + \|\partial_x^j z\|_{L_x^{\infty}(0,L;H^{\frac{1-j}{3}}(0,T))} \le C(\|w_0\|_{L^2(0,T)} + \|f\|_{L^1(0,T;L^2(0,L))})$$
 (6.6)

for j = 0, 1, 2. To this end, note that u and z can be written as

$$u(t) = W_R(t)\tilde{w}_0 - W_{bdr}(t)\vec{p}, \qquad z(t) = \int_0^t W_R(t-\tau)\tilde{f}(\tau)d\tau - W_{bdr}(t)\vec{q},$$

respectively. Here

(i)  $\tilde{w}_0$  and  $\tilde{f}$  are zero extensions of  $w_0$  and f from (0, L) to  $\mathbb{R}$ :

$$\tilde{w}_0(x) = \begin{cases} w_0(x) & x \in (0, L), \\ 0 & x \notin (0, L) \end{cases} \quad \tilde{f}_0(x, t) = \begin{cases} f(x, t) & (x, t) \in (0, L) \times (0, T), \\ 0 & x \notin (0, L). \end{cases}$$

(ii)  $W_R(t)$  is the  $C_0$  semigroup associated to the initial value problem

$$\mu_t + \mu_{xxx} = 0$$
,  $\mu(x,0) = \tilde{w}_0(x)$ ,  $x \in R$ ,  $t \in (0,T)$ .

(iii)  $\vec{p} = (p_1, p_2, p_3)$  with

$$p_1(t) = \mu_{xx}(0,t), \quad p_2(t) = \mu(L,t), \quad p_3(t) = \mu_x(L,t)$$

where  $\mu(t) = W_R(t)\tilde{w}_0$ .

(iv)  $\vec{q} = (q_1, q_2, q_3)$  with

$$q_1(t) = \tilde{z}_{xx}(0,t), \quad q_2(t) = \tilde{z}(L,t), \quad q_3(t) = \tilde{z}_x(L,t)$$

where

$$\tilde{z} = \int_0^t W_R(t-\tau)\tilde{f}(\tau)d\tau.$$

According to [14], for j = 0, 1, 2,

$$\|\partial_x^j \mu\|_{L^{\infty}(\mathbb{R}; H^{\frac{1-j}{3}}(0,T))} \le C \|\tilde{w}_0\|_{L^2(\mathbb{R})} \le C \|w_0\|_{L^2(0,L)}$$

and

$$\|\partial_x^j \tilde{z}\|_{L^\infty_x(\mathbb{R}; H^{\frac{1-j}{3}}(0,T))} \le C \|\tilde{f}_0\|_{L^1(0,T;L^2(\mathbb{R}))} \le C \|f_0\|_{L^1(0,T;L^2(0,L))}.$$

Furthermore, by Lemma 6.1,

$$\|\partial_x^j W_{bdr}(t)\vec{p}\|_{L^{\infty}(0,L;H^{\frac{1-j}{3}}(0,T))} \le C \|\vec{p}\|_{\mathcal{K}_T} \le C \|w_0\|_{L^2(0,L)}$$

and

$$\|\partial_x^j W_{bdr}(t) \vec{q}\|_{L^{\infty}_x(0,L;H^{\frac{1-j}{3}}(0,T))} \le C \|\vec{q}\|_{\mathcal{K}_T} \le C \|f\|_{L^1(0,T;L^2(0,L))}.$$

The proof of Proposition 2.8 is thus complete. ■

**Proof of Lemma 4.2**. As in the above proof (see also [1]), the solution  $\mu$  of

$$\begin{cases} \mu_t + \mu_{xxx} = 0, & \mu(x,0) = 0, \ (x,t) \in (0,L) \times (0,T), \\ \mu(0,t) = 0, & \mu(L,t) = g_2(t), & \mu_x(L,t) = 0, & t \in (0,T) \end{cases}$$

can be written as

$$\mu(x,t) = \mu_1(x,t) + \mu_2(x,t) + \mu_3(x,t)$$

with

$$\mu_j(x,t) = \frac{3}{\pi} \mathcal{R}e \int_0^\infty e^{i\rho^3 t} e^{\lambda_j(\rho)x} S_j(\rho) \rho^2 \hat{g}_2(i\rho^3) d\rho$$

for j = 1, 2, 3 where

$$\lambda_1(\rho)=i\rho,\quad \lambda_2(\rho)=\frac{\sqrt{3}}{2}\rho-\frac{1}{2}i\rho,\quad \lambda_3(\rho)=-\frac{\sqrt{3}}{2}\rho-\frac{1}{2}i\rho,$$

$$S_1(\rho)\sim 1, \quad S_2(\rho)\sim e^{-\frac{\sqrt{3}}{2}L\rho}, \quad S_3(\rho)\sim 1, \quad as \ \rho\to +\infty.$$

Arguing as before,  $g_2 \in H^{\frac{1}{3}}(0,T)$  implies that  $\mu \in X_T$ . On the other hand, if  $\mu \in X_T$ , we show that we must have  $g_2 \in H^{\frac{1}{3}}(0,T)$ . First note that as

$$\mu_{1}(x,t) = \frac{3}{\pi} \mathcal{R}e \int_{0}^{\infty} e^{i\rho^{3}t} e^{i\rho x} S_{1}(\rho) \rho^{2} \hat{g}_{2}(i\rho^{3}) d\rho$$
$$= \frac{1}{\pi} \mathcal{R}e \int_{0}^{\infty} e^{i\nu t} e^{i\nu^{\frac{1}{3}}x} S_{1}(\nu^{\frac{1}{3}}) \hat{g}_{2}(i\nu) d\nu$$

and

$$\partial_x \mu_1(x,t) = \frac{1}{\pi} \mathcal{R}e \int_0^\infty e^{i\nu t} e^{i\nu \frac{1}{3}x} \nu^{\frac{1}{3}} S_1(\nu^{\frac{1}{3}}) \hat{g}_2(i\nu) d\nu,$$

it follows from the Plancherel Theorem that for a constant c > 0

$$\|\partial_x \mu_1\|_{L_x^2(0,L;L_t^2(R))}^2 = c\|g_2\|_{H^{\frac{1}{3}}(0,T)}^2.$$

Therefore,  $\mu_1 \in L^2(0,T;H^1(0,L))$  if and only if  $g_2 \in H^{\frac{1}{3}}(0,T)$ . Regarding  $\mu_2$ , as

$$\partial_x \mu_2(x,t) = \frac{1}{\pi} \mathcal{R}e \int_0^\infty e^{i\nu t} \exp\left(-\frac{\sqrt{3}}{2} \nu^{\frac{1}{3}} (L-x) - \frac{1}{2} ix \nu^{\frac{1}{3}}\right) \nu^{\frac{1}{3}} S_2(\nu^{\frac{1}{3}}) e^{\frac{\sqrt{3}}{2} \nu^{\frac{1}{3}} L} \hat{g}_2(i\nu) d\nu,$$

we obtain

$$\|\partial_x \mu_2(x,\cdot)\|_{L^2_t(\mathbb{R})}^2 = c \int_0^\infty e^{-\sqrt{3}\nu^{\frac{1}{3}}(L-x)} \nu^{\frac{2}{3}} |\hat{g}_2(i\nu)|^2 d\nu$$

and

$$\int_0^L \|\partial_x \mu_2(x,\cdot)\|_{L^2_t(\mathbb{R})}^2 dx = c \int_0^\infty \nu^{\frac{1}{3}} |\hat{g}_2(i\nu)|^2 d\nu = c \|g_2\|_{H^{\frac{1}{6}}(0,T)}.$$

Similarly, we also have

$$\int_0^L \|\partial_x \mu_3(x,\cdot)\|_{L^2_t(\mathbb{R})}^2 dx = c \int_0^\infty \nu^{\frac{1}{3}} |\hat{g}_2(i\nu)|^2 d\nu = c \|g_2\|_{H^{\frac{1}{6}}(0,T)}.$$

Hence,  $\mu_2 + \mu_3 \in X_T$  if and only if  $g_2 \in H^{\frac{1}{6}}(0,T)$ . Consequently,  $\mu \in X_T$  if and only if  $g_2 \in H^{\frac{1}{3}}(0,T)$ . The proof of Lemma 4.2 is complete.  $\blacksquare$ 

#### 6.2Proof of Lemma 3.3

If  $N_T \neq \{0\}$ , then the map  $\varphi_T \in \mathbb{C}N_T \to A(\varphi_T) \in \mathbb{C}N_T$  has at least one eigenvalue. Therefore there exist  $\lambda \in \mathbb{C}$  and  $\varphi_0 \in H^3(0,L) \setminus \{0\}$  such that

$$\begin{cases} \lambda \varphi_0 = -\varphi_0' - \varphi_0''', \\ \varphi_0(0) = 0, \ \varphi_0'(0) = 0, \ \varphi_0(L) + \varphi_0''(L) = 0, \ \varphi'(L) = 0. \end{cases}$$
(6.7)

The solution of (6.7) satisfies  $\varphi_0(x) = \sum_{j=1}^3 C_j e^{\mu_j x}$  with  $\mu_j$  the roots of the polynomial

$$P(\mu) = \lambda + \mu + \mu^3.$$

More explicitly, they satisfy

$$\begin{cases} \mu_1 + \mu_2 + \mu_3 = 0\\ \mu_1 \mu_2 + \mu_2 \mu_3 + \mu_3 \mu_1 = 1\\ \mu_1 \mu_2 \mu_3 = \lambda \end{cases}$$
(6.8)

and  $C_j$  for j = 1, 2, 3 are the solutions of the system

$$\begin{pmatrix} 1 & 1 & 1 \\ \mu_1 & \mu_2 & \mu_3 \\ (1+\mu_1^2)e^{\mu_1 L} & (1+\mu_2^2)e^{\mu_2 L} & (1+\mu_3^2)e^{\mu_3 L} \\ \mu_1 e^{\mu_1 L} & \mu_2 e^{\mu_2 L} & \mu_3 e^{\mu_3 L} \end{pmatrix} \begin{pmatrix} C_1 \\ C_2 \\ C_3 \end{pmatrix} = \begin{pmatrix} 0 \\ 0 \\ 0 \end{pmatrix}.$$

Let us denote,  $a = L\mu_1$  and  $b = L\mu_2$ . Then by (6.8),  $c = L\mu_3 = -L(a+b)$  and

$$L^2 = -(a^2 + ab + b^2).$$

Reducing the rows of the matrix, one obtains the new one

$$M := \begin{pmatrix} 1 & 1 & 1 & 1 \\ 0 & b-a & -2a-b \\ 0 & (L^2+b^2)e^b - (L^2+a^2)e^a & (L^2+(a+b)^2)e^{-(a+b)} - (L^2+a^2)e^a \\ 0 & be^b - ae^a & -(a+b)e^{-(a+b)} - ae^a, \end{pmatrix}.$$

The system has non-zero solutions if  $det(M) \neq 0$ , which implies

$$\frac{(a+b)e^{-(a+b)} + ae^a}{2a+b} = \frac{be^b - ae^a}{b-a},$$

$$\frac{abe^{-(a+b)} + b(a+b)e^a}{2a+b} = \frac{b(a+b)e^a - a(a+b)e^b}{b-a}.$$
(6.9)

$$\frac{abe^{-(a+b)} + b(a+b)e^a}{2a+b} = \frac{b(a+b)e^a - a(a+b)e^b}{b-a}.$$
 (6.10)

Simplifying (6.9) one gets that

$$e^{-(a+b)} = \frac{2a+b}{b^2-a^2}be^b - \frac{2b+a}{b^2-a^2}ae^a$$
$$= \frac{(a+b)}{(b-a)}\left(\frac{2a+b}{b}e^b - \frac{2b+a}{a}e^a\right)$$

and from (6.10), we obtain

$$(2a+b)(\frac{b}{b+a} - \frac{a+b}{b})e^b = (2b+a)(\frac{a}{b+a} - \frac{a+b}{a})e^a$$
$$e^b = \frac{b^2}{a^2}e^a.$$

Therefore, the set of non-zero solution is empty if and only if L does not belong to

$$\mathbb{F} = \left\{ L \in \mathbb{N} : L^2 = -(a^2 + ab + b^2) \text{ with } a, b \in \mathbb{C}^2 \text{ satisfying} \quad \frac{e^a}{a^2} = \frac{e^b}{b^2} = \frac{e^{-(a+b)}}{(a+b)^2} \right\},$$

which concludes the proof of Lemma 3.3.

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